# MOMENT GROWTH BOUNDS ON CONTINUOUS TIME MARKOV PROCESSES ON NON-NEGATIVE INTEGER LATTICES 

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#### Abstract

We consider time homogeneous Markov processes in continuous time with state space $\mathbb{Z}_{+}^{N}$ and provide two sufficient conditions and one necessary condition for the existence of moments $E\left(\|X(t)\|^{r}\right)$ of all orders $r \in \mathbb{N}$ for all $t \geq 0$. The sufficient conditions also guarantee an exponential in time growth bound for the moments. The class of processes studied has finitely many state independent jumpsize vectors $\nu_{1}, \ldots, \nu_{M}$. This class of processes arises naturally in many applications such as stochastic models of chemical kinetics, population dynamics and epidemiology for example. We also provide a necessary and sufficient condition for stoichiometric boundedness of species in terms of $\nu_{j}$.


1. Introduction. Time homogeneous Markov processes in continuous time with the non-negative integer lattice as state space arise in stochastic models of chemical kinetics, predator-prey systems, and epidemiology, etc. While the primary focus of this paper will be the Markov processes describing stochastic chemical kinetics, the results derived will be of use in other applications where the processes have similar structure. More specifically, any time homogeneous Markov process model that evolves in continuous time on the state space $\mathbb{Z}_{+}^{N}$ and has finitely many types of jump events with fixed (state and time independent) jump sizes $\nu_{1}, \ldots, \nu_{M}$ will be the subject of study in this paper.

A stochastic chemical system with $N \in \mathbb{N}$ species and $M \in \mathbb{N}$ reaction channels is described by a Markov process $X(t)$ in continuous time $t \geq 0$ with state space $\mathbb{Z}_{+}^{N}$. The $i$ th component $X_{i}(t)$ describes the (random) number of species at time $t$. The probability law of the process is uniquely characterized by the stoichiometric matrix $\nu$ which is $N \times M$ with integer entries and the propensity function $a: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}_{+}^{M}$. The functions $a_{j}: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}_{+}$are also known as intensity functions or rate functions. We shall use the term propensity, which is used in chemical kinetics. The function $a_{j}(x)$ describes the "probabilistic rate" at which reaction $j$ occurs while in state $x$. More precisely, given

[^0]$X(t)=x$, the probability that reaction $j$ occurs during $(t, t+h]$ is given by $a_{j}(x) h+o(h)$ as $h \rightarrow 0+$. Column vectors of $\nu$ are denoted by $\nu_{j}$ for $j=1, \ldots, M$, and $\nu_{j}$ describes the change of state due to one occurrence of reaction $j$. See [6, 9 ] for a general introduction to stochastic models in chemical kinetics.

As an example consider the system with $N=2$ species $S_{1}, S_{2}$ and $M=2$ reactions given by

$$
S_{1} \rightarrow S_{2}, \quad S_{2} \rightarrow S_{1}
$$

Here the first reaction is one where one $S_{1}$ is converted into one $S_{2}$, and the second reaction is precisely the reversal of the first. The stoichiometric vectors are given by $\nu_{1}=$ $(-1,1)^{T}$ and $\nu_{2}=(1,-1)^{T}$. In the standard model of chemical kinetics the propensity functions for this example are given by $a_{1}(x)=c_{1} x_{1}$ and $a_{2}(x)=c_{2} x_{2}$, and in general the propensity functions are derived from combinatorial considerations and hence are polynomials [6]. In this paper however we allow a more general form for the propensities, as we do not want to limit ourselves to models arising in chemical kinetics.

The time evolution of the probabilities $p(t ; x)=\operatorname{Prob}(X(t)=x)$ is governed by Kolomogorov's forward equations

$$
\begin{equation*}
\frac{d}{d t} p(t ; x)=\sum_{j=1}^{M}\left[a_{j}\left(x-\nu_{j}\right) p\left(t ; x-\nu_{j}\right)-a_{j}(x) p(t ; x)\right] \tag{1.1}
\end{equation*}
$$

where the functions $a_{j}$ are understood to be zero if $x-\nu_{j} \notin \mathbb{Z}_{+}^{N}$, and this is typically an infinite system of ODEs indexed by $x \in \mathbb{Z}_{+}^{N}$. While the initial condition in general may be an arbitrary initial distribution $p(0 ; x)$ on $\mathbb{Z}_{+}^{N}$, it is adequate to study the case of deterministic initial conditions, i.e. $p(0 ; x)=\delta_{x_{0}}(x)$, in order to make conclusions about the general case.

In many practical examples, the system is bound to stay in a finite subset of $\mathbb{Z}_{+}^{N}$ which is determined by the initial state $x_{0}$. In the above example $S_{1} \rightarrow S_{2}, S_{2} \rightarrow S_{1}$, it is clear that the total number of species $X_{1}(t)+X_{2}(t)$ is conserved for all time $t \geq 0$. As a result the system shall remain in a finite subset of $\mathbb{Z}_{+}^{N}$. While such conservation laws and the consequent boundedness of the system are easy to spot for small systems, it may be difficult to decide for a large system. In this paper we develop a systematic theory of boundedness of species and provide necessary and sufficient conditions based on results from the study of convex cones in finite dimensions. These conditions are expressed in terms of the solution of linear inequalities which can be formulated as a linear programming problem for which several solution techniques exist [10.

While (1.1) forms a linear system of equations, analytical or even numerical computations of $p(t ; x)$ are often unwieldy even for bounded systems. In applications it is often of interest to know the moments $E\left(\|X(t)\|^{r}\right)$ for $r \in \mathbb{N}$ where $\|$.$\| is some norm on \mathbb{R}^{N}$. When the propensity functions are linear (or affine), it is possible to derive evolution equations for the moments which are closed. However, for nonlinear propensities it is not straightforward to even decide if the system has finite moments, let alone compute them.

The time evolution of the expected value of some function $h$ of the state $E(h(X(t)))$ satisfies the so-called Dynkin's formula

$$
\frac{d}{d t} E(h(X(t)))=\sum_{j=1}^{M} E\left[\left(h\left(X(t)+\nu_{j}\right)-h(X(t))\right) a_{j}(X(t))\right] .
$$

While it is tempting to use $h(x)=\|x\|^{r}$ to derive the time evolution of the $r$ th moment $E\left(\|X(t)\|^{r}\right)$, care must be taken, as the above equation may not hold for unbounded functions $h$. In this paper we derive with care some sufficient conditions for the moments $E\left(\|X(t)\|^{r}\right)$ to exist for all $r \in \mathbb{N}$ and satisfy an exponential (in time) growth bound. We also provide a necessary condition for moments $E\left(\|X(t)\|^{r}\right)$ to exist for all $r$ and all $t \geq 0$.

A set of sufficient conditions under which a large class of queueing networks (which are time inhomogeneous Markov processes on $\mathbb{Z}_{+}^{N}$ ) has moments converging asymptotically as $t \rightarrow \infty$ is obtained in 3]. A recent work [2] obtains a set of sufficient conditions under which $\sup _{t \geq 0} E\left(\|X(t)\|^{r}\right)<\infty$. The results obtained in this paper are for existence of moments for all finite $t \geq 0$ without requiring that $\sup _{t \geq 0} E\left(\|X(t)\|^{r}\right)<\infty$. This allows for systems which experience exponential growth (in time). Some sufficient conditions for the existence of moments for all finite $t \geq 0$ in the form of one-sided Lipschitz condition may be found for stochastic differential equations (SDEs) driven by Brownian motion in [7,8]. The class of processes studied in this paper is of a different form, and consequently our results are of a different flavor.

The rest of the paper is organized as follows. In Section 2 we develop some mathematical preliminaries and provide necessary and sufficient conditions for what we call the stoichiometric boundedness of species. The analysis in this section is purely deterministic. In Section 3 we provide three main results, two sufficient conditions and a necessary condition for the existence of all moments for all time $t \geq 0$. We illustrate our results via examples where appropriate.
2. Preliminaries and boundedness of species. A chemical system or a system is characterized by a stoichiometric matrix $\nu \in \mathbb{Z}^{N \times M}$ and a propensity function $a: \mathbb{Z}_{+}^{N} \rightarrow$ $\mathbb{R}_{+}^{M}$. When necessary the propensity function may be extended to the domain $\mathbb{Z}^{N}$ to be zero outside $\mathbb{Z}_{+}^{N}$. Associated to a chemical system and an initial condition $x \in \mathbb{Z}_{+}^{N}$ is a Markov process $X(t)$ in continuous time with $X(0)=x$ (with probability 1) as described in the introduction. We shall assume the process $X$ to have paths that are continuous from the right with left hand limits. We assume that the process $X$ is carried by a probability space ( $\Omega, \mathcal{F}$, Prob $)$.

We shall say that a propensity function is proper if it satisfies the condition that for all $x \in \mathbb{Z}_{+}^{N}$ if $x+\nu_{j} \notin \mathbb{Z}_{+}^{N}$, then $a_{j}(x)=0$. We note that properness is necessary and sufficient to ensure that the process $X$ remains in $\mathbb{Z}_{+}^{N}$ when started in $\mathbb{Z}_{+}^{N}$. We shall say that the propensity function is regular if it satisfies the condition that for all $x \in \mathbb{Z}_{+}^{N}$ and all $j=1, \ldots, M, a_{j}(x)=0$ if and only if $x+\nu_{j} \notin \mathbb{Z}_{+}^{N}$. We observe that regularity implies properness. Throughout the rest of the paper we shall assume properness. When regularity is assumed, it will be stated explicitly.

Consider a system with $N$ species reacting through $M$ reaction channels. We define the accessible set of states $\mathcal{A}_{x} \subset \mathbb{Z}_{+}^{N}$ given an initial state $x \in \mathbb{Z}_{+}^{N}$ by the condition that $y \in \mathcal{A}_{x}$ if and only if there exists $t>0$ such that

$$
\operatorname{Prob}(X(t)=y \mid X(0)=x)>0
$$

We observe that from standard Markov chain theory [1] the above definition is unchanged if the phrase "there exists $t>0$ " is replaced by "for every $t>0$ ". Furthermore $y \in \mathcal{A}_{x}$ if and only if there exists a finite sequence $\left(j_{1}, j_{2}, \ldots, j_{n}\right)$ of indices which take values in $\{1, \ldots, M\}$ such that for $y^{(l)}$ where $l=0,1, \ldots, n$ defined by

$$
y^{(l+1)}=y^{(l)}+\nu_{l}, \quad l=0,1, \ldots, n-1,
$$

with $y^{(0)}=x$, it holds that

$$
a_{j_{l}}\left(y^{(l-1)}\right)>0, \quad l=1, \ldots, n
$$

It is convenient to define the stoichiometricaly accessible set $\mathcal{S}_{x} \subset \mathbb{Z}_{+}^{N}$ given an initial state $x \in \mathbb{Z}_{+}^{N}$ by

$$
\begin{align*}
\mathcal{S}_{x} & =\left\{y \in \mathbb{R}^{N} \mid \exists v \in \mathbb{Z}_{+}^{M} \text { such that } y=x+\nu v\right\} \cap \mathbb{R}_{+}^{N} \\
& =\left\{y \in \mathbb{Z}_{+}^{N} \mid \exists v \in \mathbb{Z}_{+}^{M} \text { such that } y=x+\nu v\right\} . \tag{2.1}
\end{align*}
$$

(The second equality follows logically.)
It is clear that $\mathcal{A}_{x} \subset \mathcal{S}_{x}$. However these sets are not always equal, as seen from Example 1.

Example 1. Consider a system with $N=M=2, \nu_{1}=(3,-2)^{T}$ and $\nu_{2}=(-2,3)^{T}$. Consider the initial state $x=(1,1)^{T}$. Under the assumption of proper propensity function, at the initial state, propensities of both reactions are zero since the firing of either of the reactions will lead to a state with negative components. Thus $\mathcal{A}_{x}=\{x\}$. However $\mathcal{S}_{x}$ contains an infinite number of elements, as choosing $k=(n, n)^{T}$ where $n$ is a positive integer results in $y=x+\nu k=(1+n, 1+n)^{T}$, which are all in $\mathcal{S}_{x}$ by definition.

For $i=1, \ldots, N$ let $\pi_{i}: \mathbb{R}^{N} \rightarrow \mathbb{R}$ be the standard projection onto the $i$ th coordinate. Then if $\pi_{i}\left(\mathcal{A}_{x}\right)$ is bounded above we may conclude that species $i$ is bounded for initial condition $x$. Deciding whether $\pi_{i}\left(\mathcal{A}_{x}\right)$ is bounded above is harder than deciding whether $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above, so we shall focus on the latter first. We shall use the terminology that species $i$ is stoichiometrically bounded for the initial condition $x \in \mathbb{Z}_{+}^{N}$ provided $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above. As we shall see it turns out that stoichiometric boundedness is independent of the initial state $x$ and hence we could drop the reference to initial state when talking about stoichiometric boundedness of a species.

In order to study the sets $\mathcal{S}_{x}$ it is instructive to consider the related sets $C_{x}$ and $C_{x}^{+}$, defined as follows. Given $x \in \mathbb{R}^{N}$ we define $C_{x}$ and $C_{x}^{+}$as follows:

$$
\begin{gather*}
C_{x}=\left\{y \in \mathbb{R}^{N} \mid \exists v \in \mathbb{R}_{+}^{M} \text { such that } y=x+\nu v\right\},  \tag{2.2}\\
C_{x}^{+}=C_{x} \cap \mathbb{R}_{+}^{N} \tag{2.3}
\end{gather*}
$$

We note that $C_{x}$ is a closed convex cone with vertex $x$ and $C_{x}^{+}$is a closed convex set.
Lemma 2.1. Let $A \in \mathbb{R}^{k \times n}$ and $B \in \mathbb{Z}^{k \times n}$. Suppose there exists $v \in \mathbb{R}_{+}^{n}$ such that $A v>0$ and $B v=0$. Then there exists $w \in \mathbb{Z}_{+}^{n}$ such that $A w>0$ and $B w=0$.

Proof. Define

$$
P=\left\{u \in \mathbb{R}_{+}^{n} \mid A u>0, B u=0\right\}
$$

Note that $P$ is non-empty, a cone with vertex 0 and is relatively open in $\operatorname{ker}(B)$. Since $B$ has integer entrees and $\operatorname{ker}(B) \cap \mathbb{R}_{+}^{n}$ is non-empty, it follows that $\operatorname{ker}(B) \cap \mathbb{Q}_{+}^{n}$ is nonempty. As a relatively open set in $\operatorname{ker}(B) \cap \mathbb{R}_{+}^{n}$, the set $P$ contains elements from $\mathbb{Q}_{+}^{n}$. Since $P$ is a cone with vertex 0 , by taking a suitable positive integer multiple we can conclude $P$ contains elements from $\mathbb{Z}_{+}^{n}$.
Corollary 2.2. Let $A \in \mathbb{R}^{k \times n}$ and $B \in \mathbb{Z}^{k \times n}$. Suppose there exists $v \in \mathbb{R}_{+}^{n}$ such that $A v>0$ and $B v \geq 0$. Then there exists $w \in \mathbb{Z}_{+}^{n}$ such that $A w>0$ and $B w \geq 0$.

Proof. This follows from Lemma 2.1.
Lemma 2.3. Let $1 \leq i \leq N$ and $x \in \mathbb{R}_{+}^{N}$. Then $\pi_{i}\left(C_{x}^{+}\right)$is bounded above if and only if for every $z \in C_{x}^{+}$if $z_{i}>x_{i}$, then $z-x$ has at least one negative component.

Proof. Only if: Let $z \in C_{x}^{+}$and suppose $z_{i}>x_{i}$. If for all $\lambda>0$ the vector $x+\lambda(z-x)$ has no negative components, then we could infer that $\pi_{i}\left(C_{x}^{+}\right)$is unbounded. Thus there exists $\lambda>0$ such that $x+\lambda(z-x)$ has at least one negative component. This implies that $z-x$ has at least one negative component.

If: For any $z \in C_{x}^{+} \backslash\{x\}$ define $L_{z}$ by

$$
L_{z}=\left\{y \in \mathbb{R}_{+}^{N} \mid \exists \lambda \geq 0 \text { such that } y=x+\lambda(z-x)\right\},
$$

which will be a closed line segment (may be infinite). By assumption the set $\pi_{i}\left(L_{z}\right)$ is bounded above. To see this, if $z_{i} \leq x_{i}$, then $\pi_{i}\left(L_{z}\right)$ is bounded above by $x_{i}$. If $z_{i}>x_{i}$, then $L_{z}$ is a finite segment since $z-x$ has at least one negative component. Since $C_{x}^{+}$ may be partitioned into sets of the form $L_{z}$, what is left to be shown is that there exists a common upper bound $M>0$ such that $\pi_{i}\left(L_{z}\right)$ is bounded above by $M$ for all $z \in C_{x}^{+}$.

To see this define $f_{i}^{x}: C_{x}^{+} \backslash\{x\} \rightarrow \mathbb{R}$ by

$$
f_{i}^{x}(z)=\max \left\{y_{i} \mid y \in L_{z}\right\}=\max \left(\pi_{i}\left(L_{z}\right)\right),
$$

which is well defined. It is not difficult to show $f_{i}^{x}$ is continuous on $C_{x}^{+} \backslash\{x\}$ and constant on the sets $L_{z}, z \in C_{x}^{+} \backslash\{x\}$. Hence on the compact set

$$
\left\{y \in \mathbb{R}^{N} \mid\|y-x\|=1\right\} \cap C_{x}^{+}
$$

$f_{i}^{x}$ attains a maximum value, say $M$. It follows that $\pi_{i}\left(C_{x}^{+}\right)$has maximum value $M$.
Corollary 2.4. Let $1 \leq i \leq N$ and $x \in \mathbb{Z}_{+}^{N}$. Then $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above if and only if for every $z \in \mathcal{S}_{x}$ if $z_{i}>x_{i}$, then $z-x$ has at least one negative component.

Proof. Only if: If $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above, then so is $\pi_{i}\left(C_{x}^{+}\right)$, and by Lemma 2.3 the result follows.

If: Suppose $\pi_{i}\left(\mathcal{S}_{x}\right)$ is unbounded above. Then so is $\pi_{i}\left(C_{x}^{+}\right)$, and by Lemma 2.3 there exists $y \in C_{x}^{+}$such that $y_{i}>x_{i}$ and $y \geq x$. Hence there exists $v \in \mathbb{R}_{+}^{M}$ such that $\mu v>0$ and $\nu v \geq 0$ where $\mu$ is the $i$ th row of $\nu$. By Corollary 2.2 there exists $w \in \mathbb{Z}_{+}^{M}$ such that $\mu w>0$ and $\nu w \geq 0$. Taking $z=x+\nu w$ show that there exists $z \in \mathcal{S}_{x}$ such that $z_{i}>x_{i}$ and $z \geq x$.

Lemma 2.5. Let $1 \leq i \leq N$ and $x \in \mathbb{R}_{+}^{N}$. Then $\pi_{i}\left(C_{x}^{+}\right)$is bounded above if and only if $\pi_{i}\left(C_{0}^{+}\right)$is bounded above.

Proof. Only if: We note that $C_{x}=\{x\}+C_{0}$ and hence

$$
\left(C_{0} \cap \mathbb{R}_{+}^{N}\right)+\{x\}=C_{x} \cap\left(\mathbb{R}_{+}^{N}+\{x\}\right) \subset C_{x} \cap \mathbb{R}_{+}^{N}
$$

Thus $\pi_{i}\left(C_{0}^{+}+\{x\}\right)$ is bounded above and hence so is $\pi_{i}\left(C_{0}^{+}\right)$.
If: Suppose $\pi_{i}\left(C_{x}^{+}\right)$is unbounded above. Then by Lemma 2.3 there exists $z \in C_{x}^{+}$ such that $z_{i}>x_{i}$ and $z-x \geq 0$. This implies $z-x \in C_{0}^{+}$and $(z-x)_{i}>0$, which in turn implies that $\pi_{i}\left(C_{0}^{+}\right)$is unbounded above.

For $x \in \mathbb{Z}_{+}^{N}$ the study of $\mathcal{S}_{x}$ reduces to the study of $C_{x}^{+}$because of the following lemma.

Lemma 2.6. Let $1 \leq i \leq N$ and $x \in \mathbb{Z}_{+}^{N}$. Then $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above if and only if $\pi_{i}\left(C_{x}^{+}\right)$is bounded above.

Proof. If: This follows since $\mathcal{S}_{x} \subset C_{x}^{+}$.
Only if: Let $\mu \in \mathbb{Z}^{M}$ be such that $\mu^{T}$ is the $i$ th row of $\nu$. If $\pi_{i}\left(C_{x}^{+}\right)$is unbounded above, then by Lemma 2.3 there exists $z \in C_{x}^{+}$such that $z_{i}>x_{i}$ and $z \geq x$. Hence there exists $v \in \mathbb{R}_{+}^{M}$ such that $\mu^{T} v>0$ and $\nu v \geq 0$. From Lemma 2.1 we may conclude that there exists $w \in \mathbb{Z}_{+}^{M}$ such that $\mu^{T} w>0$ and $\nu w \geq 0$. It follows that the sequence $y^{(n)}$ defined by $y^{(n)}=x_{i}+n \mu^{T} w$ is a sequence in $\pi_{i}\left(\mathcal{S}_{x}\right)$ that tends to $+\infty$.

Corollary 2.7. Let $1 \leq i \leq N$ and $x \in \mathbb{Z}_{+}^{N}$. Then $\pi_{i}\left(\mathcal{S}_{x}\right)$ is bounded above if and only if $\pi_{i}\left(C_{0}^{+}\right)$is bounded above.

Proof. This follows from Lemmas 2.5 and 2.6 .
Suppose a certain non-negative linear combination

$$
\alpha_{1} X_{1}(t)+\cdots+\alpha_{N} X_{N}(t)
$$

of species is always non-increasing with time and suppose $\alpha_{i}>0$. Then we can write

$$
X_{i}(t) \leq\left(1 / \alpha_{i}\right) \sum_{j \neq i} \alpha_{j} X_{j}(0)
$$

to conclude that species $i$ is bounded. The existence of a non-decreasing non-negative linear combination can be equivalently stated as the existence of $\alpha \geq 0$ such that $\alpha^{T} \nu \leq 0$.

However the fact that the converse is also true is not obvious and requires results from the study of convex and cone sets as seen in the following theorem, which provides a necessary and sufficient condition for stoichiometric boundedness of a species.

Theorem 2.8. Species $i$ is stoichiometrically bounded if and only if there exists a vector $\alpha \in \mathbb{Z}_{+}^{N}$ such that $\alpha \geq 0, \alpha_{i}>0$ and $\alpha^{T} \nu \leq 0$.

Proof. We shall use Corollary 2.7 to work with $C_{0}^{+}$.
If: For all $y \in C_{0}$ there exists $v \in \mathbb{R}^{M}$ such that $v \geq 0$ and $y=\nu v$. Now suppose $y \in C_{0}^{+}$. Then $\alpha^{T} y \geq 0$. However, since $\alpha^{T} \nu \leq 0$ and $v \geq 0$, we have that

$$
\alpha^{T} y=\alpha^{T} \nu v \leq 0
$$

Thus $\alpha^{T} y=0$. Since $\alpha_{i}>0$ it follows that $y_{i}=0$. Hence $\pi_{i}\left(C_{0}^{+}\right)=\{0\}$ and is bounded above.

Only if: Define the set $B_{0}$ by

$$
B_{0}=\left\{\alpha \in \mathbb{R}^{N} \mid \alpha^{T} \nu \leq 0\right\},
$$

and note that $B_{0}=\left(C_{0}\right)^{o}$, i.e. the polar of $C_{0}$. (See the Appendix for some basics on convex analysis and definitions.) To see this, suppose $\alpha \in B_{0}$. If $y \in C_{0}$, then there exists $v \geq 0$ such that $y=\nu v$ and hence $\alpha^{T} y=\alpha \nu v \leq 0$. Hence $\alpha \in\left(C_{0}\right)^{o}$. Conversely if $\alpha \in\left(C_{0}\right)^{o}$, then for all $y \in C_{0}$ it holds that $\alpha^{T} y \leq 0$. Since $\nu_{1}, \ldots, \nu_{M} \in C_{0}$ it follows that $\alpha^{T} \nu \leq 0$ and thus $\alpha \in B_{0}$.

Since $\pi_{i}\left(C_{0}^{+}\right)$is bounded above it follows that $y_{i}=0$ for all $y \in C_{0}^{+}$. To see this, suppose $y \in C_{0}^{+}$and $y_{i}>0$. Since $C_{0}^{+}$is a cone, by taking positive multiples of $y$ we can obtain arbitrarily large elements in $\pi_{i}\left(C_{0}^{+}\right)$violating the assumption that $\pi_{i}\left(C_{0}^{+}\right)$is bounded above. Hence $e_{i} \in \mathbb{R}^{N}$ ( $i$ th standard basis vector) satisfies $e_{i}^{T} y=0 \leq 0$ for all $y \in C_{0}^{+}$and therefore by definition $e_{i} \in\left(C_{0}^{+}\right)^{o}$.

Now using $\left(\mathbb{R}_{+}^{N}\right)^{o}=-\mathbb{R}_{+}^{N}, B_{0}=\left(C_{0}\right)^{o}$ and Lemma A.1 we have that

$$
\left(C_{0}^{+}\right)^{o}=\left(C_{0} \cap \mathbb{R}_{+}^{N}\right)^{o}=\operatorname{cl}\left(\left(C_{0}\right)^{o}+\left(\mathbb{R}_{+}^{N}\right)^{o}\right)=\operatorname{cl}\left(B_{0}-\mathbb{R}_{+}^{N}\right)
$$

Since $B_{0}$ and $\mathbb{R}_{+}^{N}$ are polyhedral so is $B_{0}-\mathbb{R}_{+}^{N}$ and hence $B_{0}-\mathbb{R}_{+}^{N}$ is closed. Thus $e_{i} \in B_{0}-\mathbb{R}_{+}^{N}$. So $e_{i}=\alpha-u$ for some $\alpha \in B_{0}$ and $u \in \mathbb{R}_{+}^{N}$. Hence $\alpha=u+e_{i}$ and thus $\alpha \geq 0$ and $\alpha_{i}>0$. Since $\alpha \in B_{0}$ it follows that $\alpha^{T} \nu \leq 0$.

Thus we have shown that there exists $\alpha \in \mathbb{R}_{+}^{N}$ such that $\alpha_{i}>0$ and $\alpha^{T} \nu \leq 0$. By Corollary 2.2 it follows that we can choose such $\alpha \in \mathbb{Z}_{+}^{N}$.

The following corollary is immediate.
Corollary 2.9. A subset $I \subset\{1, \ldots, N\}$ of species is stoichiometrically bounded if and only if there exists a vector $\alpha \in \mathbb{Z}_{+}^{N}$ such that $\alpha \geq 0, \alpha_{i}>0$ for $i \in I$, and $\alpha^{T} \nu \leq 0$.

In order to facilitate the discussion of boundedness of species we shall define the notion of a counting sequence as follows. A finite sequence $\left(u_{1}, \ldots, u_{m}\right)$ where $u_{j} \in \mathbb{Z}_{+}^{M}$ is said to be a counting sequence provided $u_{1}=0$ and for $j=1, \ldots, m-1, u_{j+1}-u_{j}$ has precisely one component of value 1 with all other components being zero. The following lemma is immediate.
Lemma 2.10. Suppose the propensity function is regular. Then for given $x \in \mathbb{Z}_{+}^{N}$, a state $y \in \mathcal{A}_{x}$ if and only if there exists a counting sequence $\left(u_{1}, \ldots, u_{m}\right)$ in $\mathbb{Z}_{+}^{M}$ such that

$$
x+\nu u_{j} \geq 0, \quad j=1, \ldots, m
$$

and $y=x+\nu u_{m}$.
Proof. If $y \in \mathcal{A}_{x}$ there must be a sequence of reaction events which can move the state from $x$ to $y$ without leaving $\mathbb{Z}_{+}^{N}$. Conversely if there is such a sequence, then under the assumption of regularity of the propensity function, such a sequence will have non-zero probability of happening.

Finally we have the following theorem which relates boundedness of a species with its stoichiometric boundedness.

Theorem 2.11. (1) If species $i$ is stoichiometrically bounded, then it is bounded.
(2) Conversely if species $i$ is stoichiometrically unbounded and the propensity function is regular, then the species $i$ is unbounded for all sufficiently large initial conditions.

Proof. The first part is obvious. We shall prove the second part. Since the species $i$ is stoichiometrically unbounded, from Corollary 2.4 there exists $v \in \mathbb{Z}_{+}^{M}$ such that $\mu^{T} v>0$ and $\nu v \geq 0$ where $\mu^{T}$ is the $i$ th row of $\nu$. Let $u_{1}, u_{2}, \ldots, u_{m}$ be a counting sequence with $u_{m}=v$ and define $\bar{x} \in \mathbb{Z}_{+}^{N}$ by the condition that for $i=1, \ldots, N$,

$$
\bar{x}_{i}=\max \left\{0,-\left(\nu u_{1}\right)_{i},-\left(\nu u_{2}\right)_{i}, \ldots,-\left(\nu u_{m}\right)_{i}\right\} .
$$

Then for all $x \in \mathbb{Z}_{+}^{N}$ satisfying $x \geq \bar{x}$ we have that

$$
x+\nu u_{j} \geq 0, \quad j=1, \ldots, m .
$$

Thus it follows by Lemma 2.10 that $x+\nu v \in \mathcal{A}_{x}$. Define the sequence $\left(y^{(n)}\right)$ for $n \in \mathbb{N}$ by $y^{(n)}=x+n \nu v$. It is easy to show using mathematical induction that $y^{(n)} \in \mathcal{A}_{x}$ for all $n$ and that $y_{i}^{(n)}$ is strictly increasing so that $\pi_{i}\left(\mathcal{A}_{x}\right)$ is unbounded above.
3. Moment growth bounds. In order to facilitate the development of results concerning moment growth bounds we shall define critical species and critical reactions as follows.

We say that species $i$ is a critical species if and only if is not stoichiometrically bounded. Without loss of generality we assume that the species are ordered such that the copy number vector $x=(y, z) \in \mathbb{Z}_{+}^{N_{c}} \times \mathbb{Z}_{+}^{N-N_{c}}$, where $y$ is the copy number vector of critical species, $z$ is the copy number vector of non-critical species and $N_{c}$ is the number of critical species. A reaction channel $j$ is non-critical if and only if there exists $H: \mathbb{Z}_{+}^{N-N_{c}} \rightarrow \mathbb{R}$ such that

$$
a_{j}(x) \leq H(z)(\|y\|+1), \quad \forall x=(y, z) \in \mathbb{Z}_{+}^{N} .
$$

In other words critical reactions are those whose propensities grow faster than linearly in the critical species. We shall use $M_{c}$ to denote the number of critical reactions. Without loss of generality we shall assume that the reaction channels are ordered so that $j=1, \ldots, M_{c}$ correspond to the critical reactions.

In what follows, given a system with stoichiometric matrix $\nu$, we define the $N_{c} \times M$ matrix $\nu^{1}$, termed the critical species stoichiometric matrix, to be the submatrix of $\nu$ consisting of the rows $1, \ldots, N_{c}$ corresponding to the critical species, and we define the $\left(N-N_{c}\right) \times M$ matrix $\nu^{2}$, termed the non-critical species stoichiometric matrix, to be the submatrix of $\nu$ which consists of rows $N_{c}+1, \ldots, N$ corresponding to non-critical species. We also define the $N_{c} \times M_{c}$ matrix $\nu^{c}$, termed the critical stoichiometric matrix, to be the submatrix of $\nu$ consisting of the rows $1, \ldots, N_{c}$ corresponding to the critical species and columns $1, \ldots, M_{c}$ corresponding to critical reactions.

We first state a lemma.
Lemma 3.1. Suppose a system has regular propensity functions. Let $J \subset\{1, \ldots, M\}$ be a subset of reactions. Then the following are equivalent:
(1) There exists a norm $\|\cdot\|$ in $\mathbb{R}^{N}$ such that the following holds: for all $x \in \mathbb{Z}_{+}^{N}$ and for all $j \in J$,

$$
x+\nu_{j} \in \mathbb{Z}_{+}^{N} \Rightarrow\left\|x+\nu_{j}\right\| \leq\|x\| .
$$

(2) For the system consisting only of the reactions in $J$ all the species are stoichiometrically bounded.
(3) There exists $\alpha \in \mathbb{Z}_{+}^{N}$ such that $\alpha>0$ and $\alpha^{T} \nu_{j} \leq 0$ for all $j \in J$.

Proof. First we note that conditions 2 and 3 are equivalent by Corollary 2.9. It is also clear that 1 implies 2 (and hence 3 ). Thus it suffices to show 3 implies 1 . Suppose 3 holds. Define the norm $\|$.$\| on \mathbb{R}^{N}$ by

$$
\|x\|=\sum_{i=1}^{N} \alpha_{i}\left|x_{i}\right|
$$

Let $x \in \mathbb{Z}_{+}^{N}$ and suppose $x+\nu_{j} \in \mathbb{Z}_{+}^{N}$ for some $j \in J$. Then

$$
\left\|x+\nu_{j}\right\|=\sum_{i=1}^{N} \alpha_{i}\left(x_{i}+\nu_{i j}\right)=\sum_{i=1}^{N} \alpha_{i} x_{i}+\sum_{i=1}^{N} \alpha_{i} \nu_{i j}=\|x\|+\alpha^{T} \nu_{j} \leq\|x\| .
$$

We remark that Lemma 3.1 is typically used with $J=\left\{1, \ldots, M_{c}\right\}$, the set of critical reactions.

In order to discuss how moments $E\left(\|X(t)\|^{r}\right)$ evolve in time, first we note that the generator $\mathcal{A}$ of the Markov process with stoichiometric matrix $\nu$ and propensity function $a$ is given by

$$
\begin{equation*}
(\mathcal{A} h)(x)=\sum_{j=1}^{M}\left(h\left(x+\nu_{j}\right)-h(x)\right) a_{j}(x), \tag{3.1}
\end{equation*}
$$

where $h: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$, and we use the convention that $h(y)=0$ if $y \notin \mathbb{Z}_{+}^{N}$. We note that $\mathcal{A}$ is regarded as an operator on the Banach space $\mathcal{L}$ of bounded functions $h: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ and that the domain of $\mathcal{A}$ is not all of $\mathcal{L}$ as $a_{j}$ are typically not bounded functions. However the collection of all functions $h$ that are constant outside a compact subset of $\mathbb{Z}_{+}^{N}$ is in the domain of $\mathcal{A}$.

It follows from standard Markov process theory that for all functions $h: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ in the domain of $\mathcal{A}$ the following formula, sometimes known as Dynkin's formula, holds for all $t \geq 0$ :

$$
\begin{equation*}
\frac{d}{d t} E(h(X(t)))=\sum_{j=1}^{M} E\left[\left(h\left(X(t)+\nu_{j}\right)-h(X(t))\right) a_{j}(X(t))\right], \tag{3.2}
\end{equation*}
$$

or equivalently in integral form:

$$
\begin{equation*}
E(h(X(t)))=E(h(X(0)))+\sum_{j=1}^{M} \int_{0}^{t} E\left[\left(h\left(X(s)+\nu_{j}\right)-h(X(s))\right) a_{j}(X(s))\right] d s \tag{3.3}
\end{equation*}
$$

We suggest [4] as a general reference.

For $r \in \mathbb{N}$ we define the class $\mathbb{P}_{r}$ to be the set of functions $f: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ characterized by the condition that $f \in \mathbb{P}_{r}$ if and only if there exists $H>0$ such that

$$
|f(x)| \leq H\left(\|x\|^{r}+1\right), \quad \forall x \in \mathbb{Z}_{+}^{N}
$$

and define $\mathbb{P}_{r}^{+}$to denote the subset of $\mathbb{P}_{r}$ consisting of non-negative functions. We also define $\mathbb{P}$ by

$$
\mathbb{P}=\cup_{r \in \mathbb{Z}_{+}} \mathbb{P}_{r}
$$

and $\mathbb{P}^{+}$to denote the subset of $\mathbb{P}$ consisting of non-negative functions. We observe that the definition of classes $\mathbb{P}_{r}, \mathbb{P}$ is independent of the choice of norm on $\mathbb{R}^{N}$. We establish a few lemmas about classes $\mathbb{P}_{r}, \mathbb{P}$ first.
Lemma 3.2. Suppose $r, s \in \mathbb{Z}_{+}$and $r \leq s$. Then there exists $H>0$ such that

$$
\|x\|^{r} \leq H\|x\|^{s}, \quad x \in \mathbb{Z}_{+}^{N}
$$

Thus $\mathbb{P}_{r} \subset \mathbb{P}_{s}$.
Proof. This clearly holds for $x=0$, and for $x$ that satisfy $\|x\| \geq 1$ it holds with $H=1$. Since the set of $x$ for which $\|x\|<1$ is finite one may find $H$ large enough for this to hold for all $x \in \mathbb{Z}_{+}^{N}$.
Lemma 3.3. The classes $\mathbb{P}_{r}, \mathbb{P}$ are vector spaces (over $\mathbb{R}$ ) and any (multivariate) polynomial belongs to class $\mathbb{P}$. Suppose $f \in \mathbb{P}_{r}, y \in \mathbb{Z}^{N}$ and $g: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ is defined by $g(x)=f(x+y)$ if $x+y \in \mathbb{Z}_{+}^{N}$, else $g(x)=0$. Then $g \in \mathbb{P}_{r}$. In other words $\mathbb{P}_{r}$ (and hence $\mathbb{P})$ are shift invariant. Finally if $f \in \mathbb{P}_{r}$ and $g \in \mathbb{P}_{s}$, then $h \in \mathbb{P}_{s+r}$ where $h=f g$.

Proof. It is trivial to see that $\mathbb{P}_{r}$ is a vector space. Given $f, g \in \mathbb{P}$, by Lemma 3.2 there exists some $r \in \mathbb{Z}_{+}$such that $f, g \in \mathbb{P}_{r}$. Hence it is clear then that $\mathbb{P}$ is a vector space as well.

In order to show that all polynomials belong to $\mathbb{P}$ it is adequate to show that all monomials $p(x)=x^{\beta}$, where $\beta=\left(\beta_{1}, \ldots, \beta_{N}\right) \in \mathbb{Z}_{+}^{N}$ and

$$
x^{\beta}=x_{1}^{\beta_{1}} \ldots x_{N}^{\beta_{N}}
$$

belong to $\mathbb{P}$. Indeed

$$
\left|x^{\beta}\right| \leq\left(\|x\|_{\infty}^{\beta_{0}}+1\right), \quad \forall x \in \mathbb{Z}_{+}^{N},
$$

where $\beta_{0}=\max \left\{\beta_{1}, \ldots, \beta_{N}\right\}$. Using equivalence of norms there exists $K$ independent of $x$ such that

$$
\left|x^{\beta}\right| \leq K\left(\|x\|^{\beta_{0}}+1\right), \quad \forall x \in \mathbb{Z}_{+}^{N}
$$

To show shift invariance it is adequate to note that for $r \in \mathbb{Z}_{+}$and $x, y \in \mathbb{Z}^{N}$,

$$
\|x+y\|^{r} \leq(\|x\|+\|y\|)^{r} \leq \sum_{l=0}^{r} \frac{r!}{l!(r-l)!}\|x\|^{l}\|y\|^{r-l} \leq K_{r}\left(\|x\|^{r}+1\right)
$$

where $K_{r}$ depends on $y, r$ and is obtained in part by Lemma 3.2. Finally if $f \in \mathbb{P}_{r}$ and $g \in \mathbb{P}_{s}$ and $h=f g$, then for some $H>0$ and some $H^{\prime}>0$ independent of $x$ we have

$$
|f(x) g(x)| \leq H\left(\|x\|^{r}+1\right)\left(\|x\|^{s}+1\right) \leq H^{\prime}\left(\|x\|^{r+s}+1\right)
$$

where we have used Lemma 3.2.

Equations (3.2) and (3.3) hold for $h$ that are constant outside a compact set. The following lemma shows that under suitable assumptions these equations hold for all $h \in \mathbb{P}$.

Lemma 3.4. Let $r \in \mathbb{N}$ and suppose that $E\left(\|X(t)\|^{r}\right)<\infty$ for all $t \geq 0$. Then for every $f \in \mathbb{P}_{r}, E(|f(X(t))|)<\infty$ for every $t \geq 0$ and $E(f(X(t)))$ is continuous in $t$ for $t \geq 0$.

Suppose in addition that the propensity functions $a_{j}$ for $j=1, \ldots, M$ all belong to class $\mathbb{P}_{s}$ where $1 \leq s \leq r$. Then for each $f \in \mathbb{P}_{r-s}, E(f(X(t)))$ is continuously differentiable in $t$ for $t \geq 0$ and (3.2), (3.3) hold with $h=f$.

Proof. Given $f \in \mathbb{P}_{r}$, the claim $E(|f(X(t))|)<\infty$ is obvious under the assumptions.
To show $E(f(X(t)))$ is continuous in $t$ we first consider $f \in \mathbb{P}_{r}^{+}$. For each $K>0$ define $f^{K}: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ by $f^{K}(x)=f(x) \wedge K$, where $a \wedge b$ denotes the minimum of $a$ and $b$. Since for each $K, f^{K}$ is constant outside a compact set, Dynkin's formula (3.2) (with $h=f^{K}$ ) holds showing $E\left(f^{K}(X(t))\right)$ to be differentiable and hence continuous in $t$. Since $f^{K} \uparrow f$ as $K \uparrow \infty$, by monotone convergence $E\left(f^{K}(X(t))\right) \uparrow E(f(X(t)))$ as $K \uparrow \infty$. Hence Dini's theorem and a standard argument show that $E(f(X(t)))$ is continuous in $t$ for $t \geq 0$. For $f \in \mathbb{P}_{r}$ the proof is completed by decomposing $f$ into its positive and negative parts, $f=f^{+}-f^{-}$. Thus we have established that for every $f \in \mathbb{P}_{r}, E(f(X(t)))$ is continuous in $t$ for $t \geq 0$.

To show the second part we consider $f \in \mathbb{P}_{r-s}^{+}$, and for $K>0$ we consider the integral equation (3.3) with $h=f^{K}$. We observe that since $a_{j} \in \mathbb{P}_{s}$, if $f \in \mathbb{P}_{r-s}$, then $\mathcal{A}\left(f^{K}\right) \in \mathbb{P}_{r}, \mathcal{A} f \in \mathbb{P}_{r}$ and $f^{K}(X(t, \omega)) \rightarrow f(X(t, \omega))$ for almost all $(t, \omega)$ as $K \uparrow \infty$ where the Lebesgue measure is used for $t \geq 0$. Next we bound $\mathcal{A}\left(f^{K}\right)$ as

$$
\left|\left(\mathcal{A} f^{K}\right)(x)\right| \leq \sum_{j=1}^{M} f^{K}\left(x+\nu_{j}\right) a_{j}(x)+\sum_{j=1}^{M} f^{K}(x) a_{j}(x) \leq g(x)
$$

where

$$
g(x)=\sum_{j=1}^{M} f\left(x+\nu_{j}\right) a_{j}(x)+\sum_{j=1}^{M} f(x) a_{j}(x),
$$

and we also observe that $g \in \mathbb{P}_{r}^{+}$. Thus $E(g(X(t)))$ is finite and continuous in $t$ and thus $\int_{0}^{t} E(g(X(s))) d s<\infty$ for each $t \geq 0$. Hence the dominated convergence theorem allows us to conclude that one could take the limit as $K \rightarrow \infty$ on both sides of (3.3) with $h=f^{K}$ to conclude that the equation holds for $h=f \in \mathbb{P}_{r-s}^{+}$with all terms being finite. This shows $E(f(X(t)))$ to be continuously differentiable in $t$ for $f \in \mathbb{P}_{r-s}^{+}$. The proof is completed for $f \in \mathbb{P}_{r-s}$ by decomposing $f$ into positive and negative parts.

Lemma 3.5. Suppose $\phi:[0, \infty) \rightarrow \mathbb{R}$ is strictly positive for all $t \geq 0$, differentiable at 0 , and suppose there exist $H>0$ and $\lambda \in \mathbb{R}$ such that for all $t \geq 0$,

$$
\phi(t) \leq H \phi(0) e^{\lambda t} .
$$

Then there exists $\mu \in \mathbb{R}$ such that for all $t \geq 0$,

$$
\phi(t) \leq \phi(0) e^{\mu t}
$$

Proof. We observe that for $t>0$,

$$
\frac{\ln (\phi(t))-\ln (\phi(0))}{t} \leq \frac{\ln (H)}{t}+\lambda .
$$

We note that the right hand side is bounded for $t \geq t_{0}$ for every $t_{0}>0$, and the left hand side is bounded for $t \in\left(0, t_{0}\right]$ for some $t_{0}>0$ since the limit as $t \rightarrow 0+$ exists and is finite by assumption. Hence the left hand side is bounded for $t \in(0, \infty)$. We set

$$
\mu=\sup _{t>0}\left\{\frac{\ln (\phi(t))-\ln (\phi(0))}{t}\right\}<\infty
$$

to obtain the result.
The following theorem provides a sufficient condition for exponential moment growth bounds.

Theorem 3.6. Let $\nu^{c}$ be defined as above and suppose propensity functions belong to class $\mathbb{P}$. Suppose further that there exists $\alpha \in \mathbb{Z}_{+}^{N_{c}}$ such that $\alpha>0$ and $\alpha^{T} \nu^{c} \leq 0$. Then for each $r \in \mathbb{N}$ there exists $\mu_{r}$ such that the following holds for all $t \geq 0$ and in any norm $\|$.$\| on \mathbb{R}^{N}$ :

$$
\begin{equation*}
E\left(\|X(t)\|^{r}\right) \leq E\left(\|X(0)\|^{r}\right) e^{\mu_{r} t}+e^{\mu_{r} t}-1 \tag{3.4}
\end{equation*}
$$

Proof. First we claim that there exists $\gamma=\left(\gamma_{1}, \gamma_{2}\right) \in \mathbb{Z}_{+}^{N}$ where $\gamma_{1} \in \mathbb{Z}_{+}^{N_{c}}$ and $\gamma_{2} \in$ $\mathbb{Z}_{+}^{N-N_{c}}$ such that $\gamma>0, \gamma_{1}^{T} \nu_{j}^{1} \leq 0$ for $j=1, \ldots, M_{c}$ and $\gamma_{2}^{T} \nu_{j}^{2} \leq 0$ for $j=1, \ldots, M$ where $\nu^{1}, \nu^{2}$ are as defined earlier. To see this we first observe that by Corollary 2.9 there exists $\beta=\left(\beta_{1}, \beta_{2}\right) \in \mathbb{Z}_{+}^{N}$ where $\beta_{1} \in \mathbb{Z}_{+}^{N_{c}}$ and $\beta_{2} \in \mathbb{Z}_{+}^{N-N_{c}}$ such that $\beta_{1}^{T} \nu^{1} \leq 0$, $\beta_{2}^{T} \nu^{2} \leq 0$ and $\beta_{2}>0$. Set $\gamma_{1}=\beta_{1}+\alpha$ and $\gamma_{2}=\beta_{2}$ to obtain the desired result.

We shall use the norm defined by $\|x\|=\gamma^{T}|x|$ (where $|x|=\left(\left|x_{1}\right|,\left|x_{2}\right|, \ldots,\left|x_{N}\right|\right)$ ), and for $x \geq 0$ we have that $\|x\|=\gamma^{T} x$. For $r \in \mathbb{N}$ and $K>0$ define $f_{r}, f_{r}^{K}: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ by

$$
f_{r}(x)=\|x\|^{r}, \quad f_{r}^{K}(x)=\|x\|^{r} \wedge K
$$

We shall define $f_{r}, f_{r}^{K}$ to be zero outside $\mathbb{Z}_{+}^{N}$. It follows that $f_{r}^{K}$ are constant outside a compact set for each $K>0$ (and hence in the domain of the generator $\mathcal{A}$ ) and $f_{r}^{K} \uparrow f_{r}$ as $K \uparrow \infty$.

We write $\left(\mathcal{A} f_{r}^{K}\right)(x)=\sum_{j=1}^{M} T_{j}$ where

$$
T_{j}=\left[f_{r}^{K}\left(x+\nu_{j}\right)-f_{r}^{K}(x)\right] a_{j}(x) .
$$

We note that $a_{j}(x)=0$ and hence $T_{j}=0$ unless $x \in \mathbb{Z}_{+}^{N}$ and $x+\nu_{j} \in \mathbb{Z}_{+}^{N}$. Since we seek a non-negative upper bound for $T_{j}$ we shall consider only the case when $x \in \mathbb{Z}_{+}^{N}$ and $x+\nu_{j} \in \mathbb{Z}_{+}^{N}$.

When $j=1, \ldots, M_{c}$, due to the choice of our norm, we obtain that for $x=(y, z)$,

$$
\left\|x+\nu_{j}\right\|=\gamma_{1}^{T} y+\gamma_{2}^{T} z+\gamma_{1}^{T} \nu_{j}^{1}+\gamma_{2}^{T} \nu_{j}^{2} \leq \gamma_{1}^{T} y+\gamma_{2}^{T} z=\|x\| .
$$

Given this, we obtain that $T_{j} \leq 0$ regardless of the value of $K$.
To bound $T_{j}$ for $j=M_{c}+1, \ldots, M$, we consider the ordering of the three terms $\left\|x+\nu_{j}\right\|^{r},\|x\|^{r}$ and $K$. If $\|x\|^{r}>K$, then regardless of the value of $\left\|x+\nu_{j}\right\|$ we obtain that

$$
T_{j} \leq 0
$$

If $\|x\|^{r} \leq K$, then regardless of the value of $\left\|x+\nu_{j}\right\|^{r}$ we obtain that

$$
\begin{aligned}
T_{j} & \leq\left[\left\|x+\nu_{j}\right\|^{r}-\|x\|^{r}\right] a_{j}(x)=\left[\left(\|x\|+\gamma_{1}^{T} \nu_{j}^{1}+\gamma_{2}^{T} \nu_{j}^{2}\right)^{r}-\|x\|^{r}\right] a_{j}(x) \\
& \leq\left[\left(\|x\|+\gamma_{1}^{T} \nu_{j}^{1}\right)^{r}-\|x\|^{r}\right] a_{j}(x) \leq H(z)\left(\sum_{l=0}^{r-1} \frac{r!}{l!(r-l)!}\|x\|^{l}\left(\gamma_{1}^{T} \nu_{j}^{1}\right)^{r-l}\right)(\|y\|+1) \\
& \leq H(z) \lambda_{r}^{\prime}\left(\|x\|^{r}+1\right)=H(z) \lambda_{r}^{\prime}\left(\|x\|^{r} \wedge K+1\right),
\end{aligned}
$$

where $\lambda_{r}^{\prime}$ is a constant that does not depend on $x$ or $K$ and Lemma 3.2 has been used. On account of positivity of the above upper bound, it provides an upper bound for $T_{j}$ when $j=M_{c}+1, \ldots, M$ regardless of whether $\|x\|^{r} \leq K$ or not.

Thus we obtain the bound

$$
\left(\mathcal{A} f_{r}^{K}\right)(x) \leq\left(M-M_{c}\right) \lambda_{r}^{\prime} H(z) f_{r}^{K}(x)+\left(M-M_{c}\right) \lambda_{r}^{\prime} H(z) .
$$

Hence we obtain

$$
\frac{d}{d t} E\left(f_{r}^{K}(X(t))\right) \leq\left(M-M_{c}\right) \lambda_{r}^{\prime} E\left[H(Z(t)) f_{r}^{K}(X(t))\right]+\left(M-M_{c}\right) \lambda_{r}^{\prime} E[H(Z(t))]
$$

Using the fact that the vector copy number $Z(t)$ of the non-critical species is bounded, we obtain that

$$
\frac{d}{d t} E\left(f_{r}^{K}(X(t))\right) \leq \lambda_{r} E\left(f_{r}^{K}(X(t))\right)+\lambda_{r}
$$

where $\lambda_{r}$ is another constant. The Gronwall Lemma yields that

$$
E\left(f_{r}^{K}(X(t))\right) \leq E\left(f_{r}^{K}(X(0))\right) e^{\lambda_{r} t}+e^{\lambda_{r} t}-1, \quad t \geq 0
$$

Taking limit as $K \uparrow \infty$ and using the monotone convergence theorem we obtain

$$
E\left(f_{r}(X(t))\right) \leq E\left(f_{r}(X(0))\right) e^{\lambda_{r} t}+e^{\lambda_{r} t}-1, \quad t \geq 0
$$

Hence in the specific norm $\|x\|=\gamma^{T}|x|$ we obtain

$$
E\left(\|X(t)\|^{r}\right) \leq E\left(\|X(0)\|^{r}\right) e^{\lambda_{r} t}+e^{\lambda_{r} t}-1, \quad t \geq 0 .
$$

Using the equivalence of norms in $\mathbb{R}^{N}$ we obtain the bound (in any given norm)

$$
E\left(\|X(t)\|^{r}\right) \leq L_{r} E\left(\|X(0)\|^{r}\right) e^{\lambda_{r} t}+L_{r} e^{\lambda_{r} t}-L_{r}, \quad t \geq 0,
$$

where $L_{r}$ is a constant that depends only on the norm used and on $r$, and considering $t=0$ it is clear that $L_{r} \geq 1$. Define $\phi(t)$ by

$$
\phi(t)=E\left(\|X(t)\|^{r}\right)+1, \quad t \geq 0
$$

Then $\phi(t)>0$ for $t \geq 0$ and it follows that

$$
\phi(t) \leq L_{r} \phi(0) e^{\lambda_{r} t}+\left(1-L_{r}\right) \leq L_{r} \phi(0) e^{\lambda_{r} t} .
$$

By Lemma 3.4 it is clear that $\phi$ is continuously differentiable in $t$ for $t \geq 0$. Lemma 3.5 clinches the desired result.

Example 2. Consider the system with two species and two reactions given by $\nu_{1}=$ $(2,-1)^{T}, \nu_{2}=(-1,1)^{T}$ and $a_{1}(x)=x_{2}^{2}, a_{2}(x)=x_{1}$. Since $2 \nu_{1}+3 \nu_{2}=(1,1)^{T}$ it is clear that both species are critical (as they are stoichiometrically unbounded). However only reaction 1 is critical. Thus the critical stoichiometric matrix is the column vector $\nu_{1}$. The choice of $\gamma=(1,3)^{T}$ satisfies $\gamma^{T} \nu_{1}=-1<0$, and hence we can conclude that the moments of all orders exist and satisfy the exponential in time growth bound.

The conditions of Theorem 3.6 are not necessary to ensure that a moment growth bound of the form (3.4) holds.

Example 3. Consider a birth/death process with birth rate $a_{1}(x)=x^{m}$ and death rate $a_{2}(x)=2 x^{m}$. Then $\nu_{1}=1$ and $\nu_{2}=-1$ and the critical matrix $\nu^{c}=(1,-1)$. The conditions of Theorem 3.6 are not met if $m>1$. Nevertheless, intuitively one expects the birth rate to be compensated by the death rate of the same form but of a dominant magnitude. If we set $f_{r}(x)=x^{r}$, then

$$
\begin{aligned}
\left(\mathcal{A} f_{r}\right)(x) & =\left((x+1)^{r}-x^{r}\right) x^{m}+2\left((x-1)^{r}-x^{r}\right) x^{m} \\
& =\sum_{l=0}^{r-1} \frac{r!}{l!(r-l)!} x^{l+m}+2 \sum_{l=0}^{r-1} \frac{r!}{l!(r-l)!}(-1)^{r-l} x^{l+m} \\
& =\left(-r x^{r+m-1}+\frac{3 r(r-1)}{2} x^{r+m-2}+\ldots\right)
\end{aligned}
$$

When $m=2$ (quadratic birth/death rates) the positive term with highest power of $x$ is $3 r(r-1) x^{r} / 2$, and suitable truncation and the Gronwall Lemma may be used to obtain an exponential growth bound on all moments. If $m>2$, then the positive term $3 r(r-1) x^{r+m-2} / 2$ is a higher power than $x^{r}$ and unless $r=1$ (in which case finiteness can be shown easily regardless of $m$ ) this approach does not work.

Thus the intuition suggested above may be valid only if the propensities are quadratic at most.

Example 4. Consider a two species ( $S_{1}$ and $S_{2}$ ) model where when one $S_{1}$ and one $S_{2}$ come together one of three things can happen: the birth of an $S_{1}$, the birth of an $S_{2}$ or the death of both $S_{1}$ and $S_{2}$. This may be depicted by

$$
S_{1}+S_{2} \rightarrow 2 S_{1}+S_{2}, \quad S_{1}+S_{2} \rightarrow S_{1}+2 S_{2}, \quad S_{1}+S_{2} \rightarrow 0
$$

Thus we have $\nu_{1}=(1,0)^{T}, \nu_{2}=(0,1)^{T}$ and $\nu_{3}=(-1,-1)^{T}$. It is easy to see that both species are critical.

Suppose the propensities $a_{1}, a_{2}, a_{3}$ for these reactions are given by

$$
a_{1}(x)=a_{2}(x)=x_{1} x_{2}, \quad a_{3}(x)=2 x_{1} x_{2} .
$$

Hence all three reactions are critical. Here again the conditions of Theorem 3.6 are not met. However the fact that one occurrence of the third reaction undoes one occurrence of both of the other two and the dominant rate of the third reaction might suggest the possibility of bounded moments.

Let us use the 1-norm and set $f_{r}(x)=\|x\|^{r}$. We obtain that

$$
\left(\mathcal{A} f_{r}\right)(x)=\left[(y+1)^{r}-y^{r}\right] a_{1}(x)+\left[(y+1)^{r}-y^{r}\right] a_{2}(x)+\left[(y-2)^{r}-y^{r}\right] a_{3}(x),
$$

where $y=x_{1}+x_{2}$ and we suppose $x=\left(x_{1}, x_{2}\right) \geq 0$. Setting $a_{1}=a_{2}=a$ and $a_{3}=2 a$ and simplifying we obtain

$$
\begin{aligned}
\left(\mathcal{A} f_{r}\right)(x) & =2 a(x) \sum_{l=0}^{r-1} \frac{r!}{l!(r-l)!} y^{l}+2 a(x) \sum_{l=0}^{r-1} \frac{r!}{l!(r-l)!}(-2)^{r-l} y^{l} \\
& =\left(-r y^{r-1} a(x)+5 r(r-1) y^{r-2} a(x)+\ldots\right) .
\end{aligned}
$$

Since $a(x)=x_{1} x_{2} \leq\|x\|^{2}=y^{2}$, similarly to the $m=2$ case in Example 2, we may expect to obtain exponential growth bounds on all moments. However if $a(x)$ did not satisfy the quadratic growth bound, such bounds may not hold.

The two examples above are examples of application of the following theorem.
Theorem 3.7. Suppose that propensity functions all belong to class $\mathbb{P}$ and that there exist $C>0$ and $\gamma \in \mathbb{R}^{N}$ such that $\gamma>0$ and

$$
\gamma^{T} F(x) \leq C(\|x\|+1), \quad \forall x \in \mathbb{Z}_{+}^{N}
$$

where $F(x)=\sum_{j=1}^{M} \nu_{j} a_{j}(x)$. Further, suppose there exists $H>0$ such that for all $j$ with $\gamma^{T} \nu_{j} \neq 0$,

$$
a_{j}(x) \leq H\left(\|x\|^{2}+1\right)
$$

Then for each $r \in \mathbb{N}$ there exists $\mu_{r}$ such that equation (3.4) holds.
Proof. Define $f_{r}, f_{r}^{K}$ as in the proof of Theorem 3.6. We write $\left(\mathcal{A} f_{r}^{K}\right)(x)=\sum_{j=1}^{M} T_{j}$ where

$$
T_{j}=\left[f_{r}^{K}\left(x+\nu_{j}\right)-f_{r}^{K}(x)\right] a_{j}(x)
$$

We choose the norm defined by $\|x\|=\gamma^{T}|x|$ where $|x|=\left(\left|x_{1}\right|,\left|x_{2}\right|, \ldots,\left|x_{N}\right|\right)$. When $\|x\|^{r}>K$ we have that $T_{j} \leq 0$. When $\|x\|^{r} \leq K$, we obtain

$$
\begin{aligned}
T_{j} & \leq\left[\left(\gamma^{T} x+\gamma^{T} \nu_{j}\right)^{r}-\left(\gamma^{T} x\right)^{r}\right] a_{j}(x) \\
& =r\left(\gamma^{T} x\right)^{r-1} \gamma^{T} \nu_{j} a_{j}(x)+\sum_{l=0}^{r-2} \frac{r!}{l!(r-l)!}\left(\gamma^{T} \nu_{j}\right)^{r-l}\left(\gamma^{T} x\right)^{l} a_{j}(x) \\
& \leq r\left(\gamma^{T} x\right)^{r-1} \gamma^{T} \nu_{j} a_{j}(x)+H_{r}^{\prime}\left(\|x\|^{r}+1\right) \\
& =r\left(\gamma^{T} x\right)^{r-1} \gamma^{T} \nu_{j} a_{j}(x)+H_{r}^{\prime}\left(\|x\|^{r} \wedge K+1\right),
\end{aligned}
$$

where $H_{r}^{\prime}$ is a suitable constant. We note that for some $j$ if $\gamma^{T} \nu_{j}=0$, then there are no conditions on the form of the propensity function $a_{j}$. Otherwise the quadratic growth bound on $a_{j}$ ensures that an upper bound with highest power of at most $\|x\|^{r}$ is obtained. This leads to a bound of the form

$$
\left(\mathcal{A} f_{r}^{K}\right)(x) \leq \lambda_{r} f_{r}^{K}+\lambda_{r},
$$

where $\lambda_{r}>0$ is a suitable constant. The rest of the proof is similar to that of Theorem 3.6

The next theorem provides a necessary condition for the boundedness of all moments for all time $t \geq 0$.

Theorem 3.8. Suppose propensity functions all belong to $\mathbb{P}$ and suppose that there exist $\gamma \in \mathbb{R}^{N}, \alpha>1$ and $C>0$ such that $\gamma>0$ and

$$
\gamma^{T} F(x) \geq C\|x\|^{\alpha}, \quad \forall x \in \mathbb{Z}_{+}^{N}
$$

where $F(x)=\sum_{j=1}^{M} \nu_{j} a_{j}(x)$. Further suppose that $0 \in \mathbb{Z}_{+}^{N}$ is not both the initial and an absorbing state. Then for every $r \in \mathbb{N}$ that satisfies $a_{j}(x) \leq H\left(\|x\|^{r}+1\right)$ for all $j=1, \ldots, M$ (for some $H$ independent of $x$ ) there exists $t>0$ such that $E\left(\|X(t)\|^{r}\right)=\infty$. (As always we assume deterministic initial condition $x_{0}$.)

Proof. We shall prove by contradiction. Suppose $r \in \mathbb{N}$ satisfies $a_{j}(x) \leq H\left(\|x\|^{r}+1\right)$ for all $j=1, \ldots, M$ (for some $H$ independent of $x$ ) and assume that for all $t \geq 0$ it holds that $E\left(\|X(t)\|^{r}\right)<\infty$. Let $a_{0}=\sum_{j=1}^{M} a_{j}$. Then $a_{0} \in \mathbb{P}_{r}, E\left(a_{0}(X(t))\right)<\infty$ for $t \geq 0$, and using Lemma 3.4 it follows that $E\left(a_{0}(X(t))\right)$ is continuous in $t$. Thus we also have that

$$
\int_{0}^{t} E\left(a_{0}(X(s))\right) d s<\infty
$$

If the number of events of type $j$ occurring during $(0, t]$ is denoted by $R_{j}(t)$, then

$$
E\left(R_{j}(t)\right)=\int_{0}^{t} E\left(a_{j}(X(s))\right) d s<\infty
$$

and hence $R_{j}(t)<\infty$ with probability 1 . In other words the process is non-explosive.
First choose a norm such that $\|x\|=\gamma^{T}|x|\left(|x|=\left(\left|x_{1}\right|, \ldots,\left|x_{N}\right|\right)\right.$. By equivalence of norms, the inequality $\gamma^{T} F(x) \geq C\|x\|^{\alpha}$ still holds with possibly a different $C$. For each $K>0$ define $M_{K}>0$ as follows:

$$
M_{K}=\sup \left\{\left\|x+\nu_{j}\right\|,\|x\| \mid j=1, \ldots, M, x \in \mathbb{Z}_{+}^{n},\|x\| \leq K\right\}
$$

Clearly as $K \rightarrow \infty$, we have that $M_{K} \rightarrow \infty$.
For each $K>0$ let us introduce the function $f^{K}: \mathbb{Z}_{+}^{N} \rightarrow \mathbb{R}$ by $f^{K}(x)=\|x\| \wedge M_{K}$, and as $K \rightarrow \infty$ we have that $f^{K}(x) \rightarrow\|x\|$. We also have that

$$
\left(\mathcal{A} f^{K}\right)(x)=\sum_{j=1}^{M} \gamma^{T} \nu_{j} a_{j}(x)=\gamma^{T} F(x)
$$

for all $x \in \mathbb{Z}_{+}^{N}$ that satisfy $\|x\| \leq K$.
We define the stopping times $\tau_{K}$ by

$$
\tau_{K}=\inf \{t \mid\|X(t)\|>K\}
$$

By the non-explosivity, we have that $\tau_{K} \rightarrow \infty$ with probability 1 as $K \rightarrow \infty$.
We observe that (3.3) holds if $t$ is replaced by a bounded stopping time (4]. Since $t \wedge \tau_{K}$ is bounded above by $t$ we have that

$$
E\left(f^{K}\left(X\left(t \wedge \tau_{K}\right)\right)\right)=f^{K}\left(x_{0}\right)+\sum_{j=1}^{M} \gamma^{T} \nu_{j} E\left(\int_{0}^{t \wedge \tau_{K}} a_{j}(X(s)) d s\right)
$$

For sufficiently large $\nu_{0}>0$, we can bound

$$
\left\|\sum_{j=1}^{M} \gamma^{T} \nu_{j} a_{j}(x)\right\| \leq \nu_{0} a_{0}(x)
$$

for all $x \in \mathbb{Z}_{+}^{N}$. Since as $K \rightarrow \infty$, we have that $X\left(t \wedge \tau_{K}\right) \rightarrow X(t)$ with probability 1 , by dominated convergence theorem, we may take limits of both sides of the above integral equation to obtain that

$$
E(\|X(t)\|)=f\left(x_{0}\right)+\sum_{j=1}^{M} \gamma^{T} \nu_{j} E\left(\int_{0}^{t} a_{j}(X(s)) d s\right) .
$$

Hence we have that

$$
\frac{d}{d t} E(\|X(t)\|)=E\left(\gamma^{T} F(X(t))\right) \geq C E\left(\|X(t)\|^{\alpha}\right) \geq C(E(\|X(t)\|))^{\alpha}
$$

(the last step uses Jensen's inequality; see [11] for instance). Let $\phi(t)=E(\|X(t)\|)$. Thus $\phi$ satisfies

$$
\frac{d}{d t} \phi(t) \geq C(\phi(t))^{\alpha}, \quad t \geq 0
$$

Under the assumption on the initial state, for every $t>0, \phi(t)>0$. Fix $t_{1}>0$ to obtain $\phi(t) \geq \phi\left(t_{1}\right)>0$ for $t \geq t_{1}$. We obtain for $t \geq t_{1}$,

$$
\frac{d}{d t}\left(\frac{-(\alpha-1)}{(\phi(t))^{\alpha-1}}\right)=\frac{\frac{d}{d t} \phi(t)}{(\phi(t))^{\alpha}} \geq C
$$

and after some manipulations we obtain

$$
(\phi(t))^{\alpha-1} \geq \frac{\alpha-1}{\frac{\alpha-1}{\left(\phi\left(t_{1}\right)\right)^{\alpha-1}}-C t},
$$

showing that $\phi(t)=E(\|X(t)\|)$ is not finite for all $t>0$, reaching a contradiction.
We like to remark that the condition stated in Theorem 3.8 has not been shown to imply the non-existence of the first order moment $E(\|X(t)\|)$ for all time $t \geq 0$.

Example 5. Let us consider a modified version of Example 2 where $\nu_{1}=(2,-1)^{T}$, $\nu_{2}=(-1,1)^{T}$ and $a_{1}(x)=x_{2}^{2}$ as before, but we set $a_{2}(x)=x_{1}^{2}$. The quadratic form for $a_{2}$ makes reaction 2 also critical. In order to satisfy the sufficient condition of Theorem 3.6 we must find $\gamma \in \mathbb{Z}_{+}^{2}$ with $\gamma>0$ and $\gamma \nu \leq 0$. This requires the conditions

$$
\gamma_{1}>0, \quad \gamma_{2}>0, \quad 2 \gamma_{1} \geq \gamma_{2}, \quad \gamma_{1} \leq \gamma_{2}
$$

which cannot be met. In this example the sufficient conditions of Theorem 3.7 also lead to the same conditions on $\gamma$ which cannot be met. On the other hand if we choose $\gamma=(2,3)^{T}$, then we obtain that

$$
\gamma^{T} F(x)=a_{1}(x)+a_{2}(x)=x_{2}^{2}+x_{1}^{2}=\|x\|^{2} .
$$

Thus the condition of Theorem 3.8 is satisfied (assuming initial condition is not 0 ), and since propensities are quadratic we can conclude that $E\left(\|X(t)\|^{2}\right)=\infty$ for some $t>0$.

Appendix A. We summarize some basics from convex analysis. A set $C \subset \mathbb{R}^{n}$ is said to be convex if for each $x, y \in C$ and each $\alpha \in(0,1)$ it holds that $\alpha x+(1-\alpha) y \in C$. A set $K \subset \mathbb{R}^{n}$ is said to be a cone with vertex $x \in \mathbb{R}^{n}$ if for each $y \in K$ and each $\lambda \geq 0$ it holds that $x+\lambda(y-x) \in K$. A convex cone is simply a set that is both convex and a cone. A cone is said to be finitely generate or polyhedral if there exists a finite set of vectors $v_{1}, \ldots, v_{m} \in \mathbb{R}^{n}$ such that $y \in K$ if and only if there exist $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}$ all greater than or equal to zero such that

$$
y=x+\sum_{j=1}^{n} \lambda_{j} v_{j}
$$

where $x$ is the vertex.

It is easy to show that polyhedral cones are closed and convex sets, i.e. closed convex cones. It may be shown that if $C, K \subset \mathbb{R}^{n}$ are polyhedral cones with vertex 0 , then so are $C \cap K$ and $C+K$ where the sum of two sets is defined by

$$
C+K=\{c+k \mid c \in C, k \in K\} .
$$

The polar of a cone $K$ with vertex 0 is the cone $K^{o}$ (with vertex 0 ) defined by

$$
K^{o}=\left\{y \in \mathbb{R}^{n} \mid \forall x \in K, y^{T} x \leq 0\right\} .
$$

Lemma A.1. Suppose $C$ and $K$ are convex cones with vertex 0 in $\mathbb{R}^{n}$. Then

$$
(C \cap K)^{o}=\operatorname{cl}\left(C^{o}+K^{o}\right)
$$

Here cl refers to the closure of a set.
Proof. See [5], Exercise 2.12.

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