Estimating the division rate of the growth-fragmentation equation with a self-similar kernel

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Abstract
We consider the growth-fragmentation equation and we address the problem of estimating the division rate from the stable size distribution of the population, which is easily measured, but non-smooth. We propose a method based on the Mellin transform for growth-fragmentation equations with self-similar kernels. We build a sequence of functions which converges to the density of the population in division, simultaneously in several weighted $L^2$ spaces, as the measurement error goes to 0. This improves the previous results for self-similar kernels and allows us to understand the partial results for general fragmentation kernels. Numerical simulations confirm the theoretical results. Moreover, our numerical method is tested on real biological data, arising from a bacteria growth and fission experiment.

Keywords: growth-fragmentation equations, self-similar kernel, regularization, Mellin transform

(Some figures may appear in colour only in the online journal)
**Introduction**

A major issue for the coming years is to demonstrate the suitability of mathematical models to represent biological phenomena in a fully quantitative way. In this respect, inverse problem (IP) methods have already proved to be efficient. As a contribution to this field, and following previous studies [11, 12, 14, 26], this paper proposes an improved method to calibrate a growth model of major importance: the growth-fragmentation equation. We also apply our method to the experimental data of bacterial growth, as a proof of concept for its accuracy.

To describe the growth and division of particles over time, one of the key equations in the field of structured population dynamics is the so-called growth-fragmentation or cell division equation. This equation was first introduced at the end of the sixties to model cells dividing by fission [4], but it is also used to model protein polymerization [6], neuron networks [23, 24] and the TCP/IP window size protocol for the internet [1]. For all these application fields, the common point is that the ‘particle’ under concern (which can be cells, polymers, dusts, windows, etc) is well-characterized by its ‘size’, i.e. a one-dimensional quantity which grows over time, and which is distributed among the offspring of the particle when it divides.

The population is described by its concentration of particles of size $x$ at time $t$, denoted by $n(t, x)$. The equation for $n$ is obtained either by a mass balance, in the same spirit as for fluid dynamics [3, 20], or by considering the Kolmogorov equation for the underlying jumping process [8, 11]:

$$
\begin{align*}
\frac{\partial}{\partial t} n(t, x) + \frac{\partial}{\partial x} (g(x) n(t, x)) &= -B(x) n(t, x) + k \int_{0}^{\infty} \kappa(x, y) B(y) n(t, y) \, dy, \quad x \geq 0, \\
n(0, x) &= n^0(x), \quad g(0)n(t, 0) = 0.
\end{align*}
$$

(1)

Particles of size $x$ grow with a growth rate $g(x)$ and divide with a division rate $B(x)$. When a division occurs for a particle of size $y$, it splits into an average of $k \geq 1$ smaller particles, giving rise to a particle of size $x$ with a probability rate $\kappa(x, y)$. This leads to the following assumptions:

$$
\kappa(x, y) = 0 \quad \text{if} \ x > y \quad \text{and} \quad \int_{0}^{\infty} \kappa(x, y) \, dx = 1.
$$

Given an initial datum $n^0 \in L^1$, the existence of a unique solution $n$ in $C([0, +\infty); L^1)$ to equation (1) follows from the classical analysis of transport equations, while additional integrability and $L^\infty$ bounds may follow from the entropy structure of such models [21].

As concerns the asymptotic behaviour, under fairly general balance assumptions on the parameters $g$, $B$ and $\kappa$, it was proved that the population grows exponentially over time but tends to a steady profile, i.e. there exists a unique $\lambda > 0$ and a unique function $N \in L^1([0, +\infty))$ solution of the following eigenvalue problem:

$$
\begin{align*}
\frac{\partial}{\partial x} (g(x)N(x)) + (B(x) + \lambda) N(x) - k \int_{0}^{\infty} \kappa(x, y) B(y) N(y) \, dy &= 0 \quad x \geq 0, \\
g(0)N(0) &= 0, \quad N(x > 0) > 0, \quad \lambda > 0, \quad \int_{0}^{\infty} N(x) \, dx = 1,
\end{align*}
$$

(2)

and for a certain weighted norm, linked to the adjoint eigenvector, we have:

$$
e^{-\lambda t} n(t, x) \to N(x),$$

see [10, 21, 25] for more details, and recently [2, 5, 19, 22–24] for proof of an exponential speed of convergence. Remarkably, this trend to equilibrium is not only a mathematical result but is also supported by biological evidence [27], which may explain the very fast desynchronization of cells, for instance [7].
For the problem that we are dealing with here, i.e. the calibration of the parameters of the equation, it is thus valid to consider that we start with a noisy measure of the steady behaviour, i.e. the exponential rate $\lambda$, also called the fitness of the population, and the steady profile $N(x)$. The way we can use a measure $\lambda$, on the fitness $\lambda$ to determine the time-scale of the equation, for instance through a constant $c$ multiplying the growth rate $g(x)$, was explained in [12, 13]: for the sake of simplicity, and in order to concentrate on the novel aspects of our study, we do not detail this aspect further, and here we consider that $\lambda > 0$ as well as $g(x)$ are fully known. We model the measurement error in a deterministic way, assuming that we are given an observation $N^\epsilon$ such that:

$$\|N - N^\epsilon\|_{H^s([0, +\infty))} \leq \epsilon, \text{ for some } s \in [-1, 0],$$

where $\| \|$ stands for the norm of a certain Sobolev Hilbert space $H^s([0, +\infty))$.

As in the previous studies, we suppose that $g(x)$, $k$ and $\kappa(x, y)$ are already known, or guessed at—such measures may be carried out directly in many cases [27]—and we concentrate on estimating the division rate $B(x)$.

We note that $B$ only appears multiplied by $N$ in equation (2), so that it cannot be accurately estimated from this equation where $N$ vanishes (near 0 and near $+\infty$). Hence, we denote $H = BN$ and in this paper we focus on estimating $H$ rather than $B$ (the interested reader may find in [12] a fully rigorous estimate of $B$, obtained after a truncated division by $N$, i.e. by defining $B_\epsilon = \frac{H}{N^{\frac{1}{k} > 3}}$). Equation (2) may be formulated in terms of $H$:

$$\frac{\partial}{\partial x} (g(x)N(x)) + \lambda N(x) = -H(x) + k \int_0^{+\infty} \kappa(x, y)H(y) \, dy, \quad x \geq 0. \tag{4}$$

We can now formulate precisely the IP under study, which the previous simplifications make linear.

**Inverse Problem (IP).** Given a measurement $N^\epsilon$ of the solution $N$ of equation (2) such that $N^\epsilon - N$ satisfies estimate (3), how can we get an approximation $H^\epsilon$ of $H$ solution of equation (4)?

In [14, 26], a theoretical and numerical solution was proposed for the specific case of equal mitosis\(^6\), i.e. when $k = 2$ and $\kappa(x, y) = \delta_{y=x}$, and estimates were obtained in $L^2([0, +\infty))$. In practice, the numerical solution exhibited oscillatory behaviour for large values of $x$. Then, a generalization was proposed in [15] for a general division kernel $\kappa$, but the estimation was proved in spaces $L^2(x^q \, dx)$ with $q > 3$, which leads in practice to a noise amplification around zero. These two problems lead us to look for an estimate in a space that would guarantee smooth behaviour both around zero and infinity: we are looking for an error estimate valid in $L^2(x^q \, dx)$ spaces, with both $q = 0$ and $q > 3$, in order to avoid noise amplification around both zero and infinity.

The aim of this paper is to give a complete solution of the IP for general self-similar fragmentation kernels. Self-similarity refers here to the cases where fragmentation only depends on the ratio between the parent size and the offspring size. Mathematically, this means that there exists a probability measure $\kappa_0$ such that:

$$\kappa(x, y) = \frac{1}{y} \kappa_0\left(\frac{x}{y}\right), \quad \kappa_0 \in P([0, 1]), \quad \int_0^1 z \, d\kappa_0(z) = \frac{1}{k}. \tag{5}$$

\(^5\) In [14, 15, 26], the norm was in $L^2$, i.e. $s = 0$. In [12], the noise coming from density estimation could be heuristically compared to $H^{-1/2}$, i.e. $s = -1/2$.

\(^6\) This kernel was the most widely studied (and is relevant for Escherichia coli or TCP/IP for instance).
Let us give a weak formulation of equation (4) when \( \kappa \) satisfies assumption (5). For every \( \varphi \) in \( C_0^\infty((0, +\infty)) \):

\[
- \int_0^{+\infty} gN\varphi' \, dx + \lambda \int_0^{+\infty} N\varphi \, dx + \int_0^{+\infty} H\varphi \, dx = k \int_0^{+\infty} H(x) \int_0^1 \varphi(xy) \, dx \kappa_0(y) \, dy.
\]

(6)

For a self-similar kernel \( \kappa(x, y) = \frac{1}{y} \kappa_0\left(\frac{x}{y}\right) \), the Mellin transform (cf section 3) is suitable for the problem of inverting the linear operator \( L \) because the term \( KH = \int_0^{+\infty} \kappa_0\left(\frac{x}{y}\right) H(y) \frac{dy}{y} \) is the multiplicative convolution of the function \( H \) and the measure \( \kappa_0 \). As the Fourier transform replaces additive convolution by product, the Mellin transform (which can be viewed as the standard Fourier transform when applying a logarithmic change of variables) replaces multiplicative convolution by product. See section 2.

In the first section, we give the main steps required for regularizing and solving the IP, and state the main result. The second and third sections are devoted to the Mellin transform: we recall some fundamental results, which we use as the key point for our solution. The fourth section then details the numerical implementation. Technical details of the proofs are given in the appendix.

1. Estimation protocol and main result

Let us briefly explain the main steps of our method. Since it was first introduced in [26] and then developed in [12, 14, 15], we only outline its main features and we let the interested reader refer to those articles. First, we introduce some notations.

**Definition 1.1** (\( L^p_q, \mathbb{H}^q, \mathbb{W}^{n,p}, \mathbb{W}^{-n,p} \)). Let \( p \geq 1, q \in \mathbb{R}, n \in \mathbb{N}^* \).

1. For \( f : \mathbb{R} \to \mathbb{R} \) we denote

\[
\|f\|_{L^p_q} := \left( \int_0^{+\infty} |f(x)|^p |x|^q \, dx \right)^{1/p}.
\]

We define \( L^p_q \) the Banach space \( \mathbb{L}^p((0, +\infty), x^q \, dx) \) equipped with the norm \( \| \cdot \|_{L^p_q} \).

2. We denote

\[
\mathbb{W}^{n,p}((1 + x^q) \, dx) := \{ f : [0, +\infty) \to \mathbb{R} | f, \ldots, f^{(n)} \in \mathbb{L}^p((1 + x^q) \, dx) \}.
\]

When equipped with the norm

\[
\|f\|_{\mathbb{W}^{n,p}((1 + x^q) \, dx)} := \sum_{k=0}^n \|f^{(k)}\|_{\mathbb{L}^p((1 + x^q) \, dx)}.
\]

\( \mathbb{W}^{n,p}((1 + x^q) \, dx) \) is a Banach space.

3. We define similarly \( \mathbb{W}^{-n,p} (x^q \, dx) \). We also define: \( \mathbb{H}^n((1 + x^q) \, dx) = \mathbb{W}^{n,2}((1 + x^q) \, dx) \) and \( \mathbb{H}^{-n} (x^q \, dx) = \mathbb{W}^{-n,2} (x^q \, dx) \).

4. We also define:

\[
\mathbb{W}^{-1,p}((1 + x^q) \, dx) := \{ f \in \mathcal{D}'(0, +\infty) : f = g + h', g, h \in \mathbb{L}^p((1 + x^q) \, dx) \}.
\]

Equipped with the norm:

\[
\|f\|_{\mathbb{W}^{-1,p}((1 + x^q) \, dx)} := \inf_{f = g + h} \left( \|g\|_{\mathbb{L}^p((1 + x^q) \, dx)} + \|h\|_{\mathbb{L}^p((1 + x^q) \, dx)} \right),
\]

\( \mathbb{W}^{-1,p}((1 + x^q) \, dx) \) is a Banach space.
Lastly we define for \( \theta \) in \([0, 1]\) the Banach space \( \mathcal{W}^{\theta,p}(1+x^q) \) by complex interpolation:

\[
\mathcal{W}^{\theta,p}(1+x^q) := [L^p((1+x^q) \, dx), \mathcal{W}^{-1,p}(1+x^q) \, dx]_\theta.
\]

Equipped with the standard norm for complex interpolation spaces (see [31] for instance), \( \mathcal{W}^{\theta,p}(1+x^q) \) is a Banach space.

The following linear operators, \( K, L \) and \( D \), were introduced in [14];

\[
K(H) := \int_0^{+\infty} \kappa(\cdot, y) H(y) \, dy,
\]

(7)

\[
L := \text{Id} - kK,
\]

(8)

\[
D(N) := -\partial_x (gN) - \lambda N.
\]

(9)

Equation (4) may be formulated in terms of \( L \) and \( D \):

\[
L(H) = D(N).
\]

(10)

To solve the problem (IP), we thus need two steps. The first step is the same as in previous studies [12, 14, 15, 26], and consists in defining a convenient approximation \( D_\alpha(N^\varepsilon) \) of \( D(N) \) from a noisy measure \( N^\varepsilon \). The second step is inverting the operator \( L \): and finding \( u \) solution to

\[
L(u) = f.
\]

(11)

This is fully developed in section 3, and it is the key point of this article.

First step: regularization and estimation of \( D(N) \)

The operator \( N \mapsto D(N) \) defined by (9) is homogeneous to a derivative operator so that it has a de-regularizing effect, whereas the operator \( L \) is shown in section 3 to be continuously invertible from \( L^p_q \) to \( L^p_q \) spaces, under quite general assumptions on \( p, q \). Hence, starting from a datum \( N^\varepsilon \) in \( \mathcal{H}^{-s}(0, +\infty) \) for \( 0 \leq s \leq 1 \), for \( g \) which has sufficient regularity we have:

\[
D(N^\varepsilon) \in \mathcal{H}^{-s-1}(0, +\infty).
\]

To obtain an error estimate for \( H - H^\star \) in an \( L^2 \) space, we need a regularization step before applying \( L^{-1} \). This regularization being from \( \mathcal{H}^{-s-1} \) to \( L^2 = \mathbb{H}^0 \) defines a degree of ill-posedness \( a = 1 + s \) [12, 16, 32].

Among the many regularization techniques—see e.g. [18, 26]—we choose here the filtering or kernel approximation, already used in [12, 14, 15] for its simplicity and accuracy. To do so, for \( n \) in \( \mathbb{N}^+ \), we introduce a mollifier \( \rho \) satisfying:

\[
\rho \in C_0^\infty(\mathbb{R}), \quad \text{supp } \rho \subset [-1, 1], \quad \int_{-1}^{1} \rho(x) \, dx = 1,
\]

\[
\int_{-1}^{1} x^k \rho(x) \, dx = 0 \quad \text{for } k = 1, \ldots, n - 1;
\]

(12)

(the fourth condition is empty for \( n = 1 \)). For \( \alpha > 0 \) we define \( \rho_\alpha(x) = \frac{1}{\alpha} \rho \left( \frac{x}{\alpha} \right) \) that provides us with the natural approximation \( D_\alpha(N^\varepsilon) \) of \( D(N) \):

\[
D_\alpha(N^\varepsilon) := \rho_\alpha * D(N^\varepsilon),
\]

(13)

which belongs to any \( L^2_q \); all the useful estimates are recalled in lemma A.1. We have the following estimate.
Proposition 1.1. Let \( n \in \mathbb{N}^* \) and \( \rho \) satisfy assumption (12). For \( 1 \leq p < \infty, 0 \leq q < \infty \), we suppose that \( gN \in W^{p+1, q}((1+x^q)\,dx) \) and \( N \in W^{q, p}((1+x^q)\,dx) \). Let \( N^e \in W^{q, p}((1+x^q)\,dx) \) with \( 0 \leq s < 1 \). Defining \( D(N) \) by (9) and, for \( \alpha > 0 \), \( D_\alpha(N^e) \) by (13) we have:
\[
\|D_\alpha(N^e) - D(N)\|_{L^p_q} \leq C(\alpha^{-s(1+\alpha)} \|g(N^e - N)\|_{H^{-s^2}(1+x^q)\,dx}) + \alpha^{-s} \|N - N\|_{H^{-s^2}(1+x^q)\,dx}
\]
where the constant \( C \) depends on \( 2^{-s} \), \( \|\rho\|_{L^1_{-1}} \), \( \|\rho''\|_{L^1_{-1}} \), \( |\lambda| \) and \( \|gN\|_{W^{p+1, q}((1+x^q)\,dx)} \), \( \|N\|_{W^{q, p}((1+x^q)\,dx)} \). If the terms \( \|N - N\|_{H^{-s^2}(1+x^q)\,dx} \) are in the order of \( \varepsilon \), the previous estimate is optimal for \( \alpha = \varepsilon^{1/(\alpha+1+s)} \), which leads to:
\[
\|D_\alpha(N^e) - D(N)\|_{L^p_q} = O(\varepsilon^{\alpha/(\alpha+1+s)}).
\]

Second and main step: inversion of \( \mathcal{L} \)

We now define a continuous inverse \( \mathcal{L}^{-1} \) for \( \mathcal{L} \), and apply it to \( D_\alpha(N^e) \) to obtain an estimate \( H^s_\alpha \) of \( H \). In the case of a general fragmentation kernel, the operator \( \mathcal{L} : L^p_q \rightarrow L^2_q \) was shown to be continuously invertible for \( q > 3 \) in [15]. This result is generalized to \( L^p_q \) spaces, for large enough values of \( q \), in appendix E. However, even if such a result provides an error estimate for \( H - H^s_\alpha \) in \( L^p_q \) spaces with \( q > 3 \), it proves to be unsatisfactory for practical purposes; the numerical methods developed in [15] from such definition of an inverse proved to blow-up around zero. To avoid such a blow-up, we would need to define an inverse in \( L^p_q \) spaces with smaller weights \( q \), like \( q = 0 \).

This was done for the case of the equal mitosis kernel \( \kappa(x, y) = \delta_{x=y} \) in [14]; the operator \( \mathcal{L} : L^2_q \rightarrow L^2_q \) was shown to be invertible for \( q > 3 \) and also for \( q < 3 \), with two distinct definitions of the inverse—one holding true for any \( q < 3 \), the other for any \( q > 3 \).

We extend this result of [14] to any self-similar kernel, by using the Mellin transform. We give here only the main ingredients, and refer to sections 2 and 3 for full details.

Definition 1.2. Let \( \mu \) be a complex valued regular measure on \((0, +\infty)\). The Mellin transform of \( \mu \) at \( s \in \mathbb{C} \) is defined by:
\[
\mathcal{M}\mu(s) = \int_0^{+\infty} x^{s-1} \, d\mu(x),
\]
provided that this expression exists. The Mellin transform of a function \( f \) at \( s \in \mathbb{C} \) is simply obtained by taking \( \mu = f(x) \, dx \) in the previous formula.

Thanks to the relation (6) and to the definition (8) we have \( \mathcal{M}\kappa_0(s) = \mathcal{M}\kappa_0(s)\mu(s) \) and the following identity holds:
\[
\mathcal{M}\kappa_0(s) = (1 - k\mathcal{M}\kappa_0(s)) \mathcal{M}\mu(s).
\]

Formally the solution \( u \) of equation (11) may be given by
\[
\mathcal{M}u(s) = \frac{\mathcal{M}f(s)}{1 - k\mathcal{M}\kappa_0(s)}.
\]

It remains to define a continuous inverse of the Mellin transform and to apply it to the right-hand side; this can be done on a vertical line of the complex plane. To do so, we need \( 1 - k\mathcal{M}\kappa_0 \) not to vanish on such a vertical line; this is expressed in assumption 1 below.

Assumption 1. \( \zeta \in \mathbb{R} \) is such that the function \(|1 - k\mathcal{M}\kappa_0|\) is bounded from below on the vertical line \( \zeta + i\mathbb{R} \) by a positive constant.
Note that by assumption (5) on \((\kappa_0, k)\): \(1 - kM\kappa_0(1) = 1 - k < 0\) and \(1 - kM\kappa_0(2) = 0\), so that assumption (1) never holds true for \(\zeta = 2\).

Using section 2, \(M\kappa_0\) is (at least) defined in the half-plane \(\{s \in \mathbb{C} \mid \text{Re} s \geq 1\}\) and satisfies \(|M\kappa_0| \leq \frac{1}{\zeta}\) in the half-plane \(\{s \in \mathbb{C} \mid \text{Re} s \geq 2\}\). Furthermore, if \(\kappa_0((0, 1)) > 0\) then \(|M\kappa_0| < \frac{1}{\zeta}\) in the half-plane \(\{s \in \mathbb{C} \mid \text{Re} s > 2\}\); assumption 1 is always satisfied for \(\zeta > 2\).

The existence of a value \(\zeta < 2\) such that assumption 1 is satisfied is a more complex question. We will see (cf remark 3.4) that, under an additional assumption on \(M\kappa_0(q)\), assumption 1 is satisfied for almost every \(\zeta < 2\).

We then have the following proposition.

**Proposition 1.2.** If \(q \in \mathbb{R}\) is such that \(\zeta = \frac{q + 1}{2}\) satisfies assumption 1, the linear operator \(\mathcal{L}\) defined by (8) is continuously invertible from \(L^2_q\) to \(L^2_{\zeta}\).

This proposition is included in proposition 3.8 of section 3.2, where an upper bound for the norm of \(\mathcal{L}^{-1}\) is also provided.

If \((\kappa_0, k)\) satisfies (5) and if \(\kappa_0((0, 1)) > 0\) then proposition 1.2 is satisfied for every \(q > 3\) (cf section 2). Under an assumption on \(1 - kM\kappa_0(\frac{q - 1}{2} + i\epsilon)\) as \(\epsilon\) tends to \(+\infty\) we will see that every weight \(q_1 < 3\), except a countable set of values, satisfies proposition 1.2 (cf remark 3.4).

In [14] (appendix A.1) it is shown that the expression of \(\mathcal{L}^{-1}_q : L^2_q \to L^2_{\zeta}\) does not depend on \(q > 3\) or on \(q < 3\). We generalize this to self-similar kernels, cf propositions 3.2 and 3.10, respectively in \(L^2_q\) and in \(L^2_{\zeta}\) spaces.

If \(D(N) \in L^2_q\) and if \(B\) is such that \(H = BN \in L^2_q\), then \(H = \mathcal{L}^{-1}_q(D(N))\) and we can define

\[
H^e_{q,a} := \mathcal{L}^{-1}_q(D_a(N^e)) .
\]

Proposition 1.2 together with proposition 1.1 immediately gives us an estimate for \(\|H^e_{q,a} - H\|_{L^2_q}\). However, as detailed in propositions 3.2 and 3.10 of section 3 below, such an inverse will not be the same in any space \(L^2_q\), due to the fact that the function \(1 - kM\kappa_0\) has a zero at least in \(s = 2\) (which corresponds to the space \(L^2_2\)). To obtain an estimate that would be true in several \(L^2_q\) spaces, for instance for both \(q_1 = 0\) and \(q_2 > 3\), we define for a given constant \(a > 0\):

\[
H^e_{q,a}(a) := H^e_{q_1,a} 1_{[0,a)} + H^e_{q_2,a} 1_{[a,\infty)}.
\]

To obtain an estimate for \(H^e_{q,a} - H\), the last ingredient lies in the following interpolation lemma.

**Lemma 1.3.** Let \(a > 0\), and \(0 \leq q_1 \leq q_2\). Let \(f \in L^p_{q_1}\) and \(g \in L^p_{q_2}\). Then the function \(h = f 1_{[0,a]} + g 1_{[a,\infty]}\) satisfies:

\[
\forall q \in [q_1, q_2], \quad \|h\|_{L^p_q} \leq a^\frac{q_1}{p} \|f\|_{L^{p}_{q_1}} + a^\frac{q_2}{p} \|g\|_{L^{p}_{q_2}} .
\]

**Proof.** We first treat the case \(a = 1\). We have for all \(x \geq 0\):

\[
|h(x)^p| x^\theta = |f(x)|^p x^\theta 1_{[0,1]} + |g(x)|^p x^\theta 1_{(1,\infty)} \leq |f(x)|^p x^\theta 1_{[0,1]} + |g(x)|^p x^\theta 1_{(1,\infty)} .
\]

Integrating this inequality leads to \(\|h\|_{L^p_{q_1}} \leq \|f\|_{L^p_{q_1}} + \|g\|_{L^p_{q_2}},\) and the result follows from the convexity inequality \((|x|^p + |y|^p)^{\frac{1}{p}} \leq |x| + |y|\). The general case \(a > 0\) is obtained by a change of variable. We set \(z_a(x) = ax\) so that \(z_a(1) = z(a)\). We apply the above inequality to the functions \(f_a, g_a, h_a\) and the result is a consequence of the fact \(\|z_a\|_{L^p_q} = a^\frac{q_1}{p} \|z\|_{L^p_{q_1}} .\)

We can now state the main result of this paper, which is a direct consequence of propositions 1.1 and 1.2 and lemma 1.3.
Theorem 1.4. Let \((k_0, k)\) be defined by (5), and \(q_1, q_2\) be real numbers such that \(0 \leq q_1 < 3 < q_2\) and satisfy the assumptions of proposition 1.2.

Let \((\lambda, N)\) with \(N\) in \(H^s((1 + x^a) \, dx)\) and \(\lambda > 0\) be the unique solution of equation (2). We assume that: \(gN \in H^s((1 + x^{\lambda}) \, dx)\) and \(BN \in L^2((1 + x^{\lambda}) \, dx)\) for some \(n \geq 1\).

Let \(0 \leq s \leq 1\) and let \(N^\nu\) be a measure of \(N\) such that \(N^\nu \in H^{-s}((1 + x^{\lambda}) \, dx)\) for \(q = q_1, q_2\) and with a measurement error given by:

\[
\|N^\nu - N\|_{H^{-s}((1 + x^{\lambda}) \, dx)} \leq \varepsilon, \quad \|g(N^\nu - N)\|_{L^2((1 + x^{\lambda}) \, dx)} \leq \varepsilon.
\]

Let \(n \in \mathbb{N}^*\) and \(\rho\) satisfy (12). For \(\alpha > 0\), we define \(D_\rho(N^\nu)\) by (13), \(H^\varepsilon_{\rho, a}\) by (15). Let \(\alpha > 0\) and define \(H^\varepsilon_{\rho, a}\) by (16).

For every weight \(q\) in \([q_1, q_2]\) and \(0 < \alpha \leq 1\), the following estimate holds:

\[
\|H^\varepsilon_{\rho, a} - H\|_{L^2} \leq C (\alpha^{-1+\sigma} + \alpha^\sigma),
\]

where the constant \(C\) depends on \(a, q, |\lambda|, q_1, q_2, \|\rho\|_{L^2([-1,1])}, \|\rho^\sigma\|_{L^1([-1,1])}\), \(\|N\|_{L^2((1 + x^{\lambda}) \, dx)}\), \(\|gN\|_{H^s((1 + x^{\lambda}) \, dx)}\), \(\|\rho\|_{L^2([-1,1])}\), \(\|\rho^\sigma\|_{L^1([-1,1])}\).

This induces an optimal error in the order of \(e^{n/(n+1+\sigma)}\), obtained for a regularization parameter \(\alpha\) in the order of \(e^{1/(n+1+\sigma)}\).

2. Recall on the Mellin transform

Notation 2.1. For a function \(f : \xi + i\mathbb{R} \to \mathbb{C}\) we denote

\[
\|f\|_{L^p(\xi + i\mathbb{R})} = \left( \int_{\mathbb{R}} |f(\xi + i\tau)|^p \, d\tau \right)^{1/p}.
\]

The construction of the Mellin transform on \(i\mathbb{R}\) can be made in the general context of the Fourier transform on locally compact abelian groups (we refer the reader to chapter 1 of [28]). Here we consider the multiplicative abelian group \(G = (0, +\infty)\) (with unit 1), equipped with the topology inherited from \(\mathbb{R}\) and with the Haar measure \(\frac{dx}{x}\) (that is the unique measure on \(G\), up to a positive multiplicative constant, which is translation-invariant). It is easy to show that the dual group \(\Gamma\) of all the characters of \(G\) with the Gelfand topology is isomorphic to \(i\mathbb{R}\) with the topology inherited from \(\mathbb{C}\) via the application \(i\mathbb{R} \to \Gamma, \quad \mu \mapsto (x \mapsto x^{-\mu})\). In this general context, the \(L^1\) and \(L^2\) theories can be built, with the same results, as the Fourier transform on the topological group \((\mathbb{R}^n, +, dx)\). The extension of the Mellin transform to a vertical line \(\xi + i\mathbb{R}\) of the complex plane \(\mathbb{C}\) is simply obtained by defining \(\mathcal{M}f(\xi + i\tau) = \mathcal{M}g(i\tau)\) with \(g(x) = x^\xi f(x)\).

In the following subsection we state some simple definitions and results concerning the Mellin transform, which will be useful for the proofs of section 3. If a proposition cannot be found in [28] or in [30], a brief proof is sketched out in appendix B. We first introduce the Mellin transform of a measure and then the Mellin transform of a function of class \(L^1, L^2\).

2.1. Mellin transform of a measure

We recalled above in definition 1.2 the Mellin transform of a measure. For the particular case of measures supported in \([0, 1]\), this definition exists in a half-plane, as expressed in the following lemma.

Lemma 2.1. Let \(\mu\) be a non-negative measure supported in \([0, 1]\). If \(\mathcal{M}\mu(a)\) exists for some real number \(a\) then \(\mathcal{M}\mu\) is continuous in the closed half-plane \(\{s \in \mathbb{C} \mid \text{Re } s \geq a\}\) and is holomorphic in the open half-plane \(\{s \in \mathbb{C} \mid \text{Re } s > a\}\).
Definition 2.2. Let $\mu$ be a non-negative measure supported in $[0, 1]$. We define the abscissa of convergence of $\mu$ as: $\text{Abs}\mu = \inf\{a \in \mathbb{R} \mid M\mu(a) \text{ is finite}\} \in \mathbb{R} \cup \{\pm \infty\}$.

If $\text{Abs}\mu$ is finite or $-\infty$ then $M\mu$ is holomorphic in the half-plane $\{s \in \mathbb{C} \mid \text{Re } s > \text{Abs}\mu\}$.

Example 2.2.

| Self-similar kernel $\kappa_0$ | $M\kappa_0(s)$ | Abs $\kappa_0$ |
|------------------------------|----------------|-------------|
| $1_{[0,1]} \, dx$ | $s^{-1}$ | 0 |
| $\alpha x^{\alpha-1} 1_{[0,1]} \, dx$ | $\frac{s-\alpha}{\sigma^{\alpha-1}}$ | $1 - \alpha$ |
| $\delta_x, \sigma \in (0, 1)$ | $\frac{\pi}{\sigma^{\alpha-1}}$ | $-\infty$ |
| $\frac{2}{q} (x^{q-1} + (1-x)^{q-1}) 1_{[0,1]} \, dx, \; \alpha > 0$ | $\frac{1}{2\pi i} \int_{q-\infty}^{q+\infty} \varphi(s) x^{-s} ds := \frac{1}{2\pi i} \int_{\mathbb{R}} \varphi(q^* + i\tau) x^{-q^*} x^{-\tau} d\tau$ | max$(1-\alpha, 0)$ |

In the last line $B(x, y) = \int_0^1 t^{x-1} (1 - t)^{y-1} \, dt$ is the Euler Beta function which is defined for $\text{Re } x, \text{Re } y > 0$ (and linked to the Euler Gamma function by: $B(x, y) = \frac{\Gamma(x)\Gamma(y)}{\Gamma(x+y)}$). The self-similar kernel $\kappa_0 = \rho_1 \delta_{\tau_1} + \rho_2 \delta_{\tau_2}$ is studied in appendix D.

Lemma 2.3. Let $\mu$ be a non-negative measure supported in $[0, 1]$. The function $(\text{Abs}\mu, +\infty) \rightarrow \mathbb{R}, \; t \mapsto M\mu(t)$

(1) is a non-increasing and log-convex function,
(2) tends to 0 as $t$ tends to $+\infty$ if $\mu(\{1\}) = 0$,
(3) is strictly decreasing if $\mu((0, 1)) > 0$.

2.2. Mellin transform of functions on $\mathbb{C}$

For $f$ in $L^1_{q^*}$, we define the Mellin transform of $f$ by definition 1.2 applied to $\mu = f(x) \, dx$, for $s \in q^* + i\mathbb{R}$ with $q^* = q + 1$:

$$Mf(s) = \int_0^{+\infty} x^{s-1} f(x) \frac{dx}{x}.$$ 

Theorem (Riemann–Lebesgue). The Mellin transform is a linear continuous map of $L^1_{q^*}$ into $C_0^\infty(q^* + i\mathbb{R}) \subset L^\infty(q^* + i\mathbb{R})$, where $q^* = q + 1$, with norm 1.

Theorem (Inversion theorem). If $f$ is in $L^1_{q^*}$ and if $\|Mf\|_{L^1(q^*+i\mathbb{R})}$ is finite, then $q^* = q + 1$, then for a.e. $x > 0$ we have:

$$f = Mq^*^{-1}(Mf),$$

where

$$Mq^*^{-1}(\varphi(x)) := \frac{1}{2\pi i} \int_{q^*+\infty}^{q^*-\infty} \varphi(s) x^{-s} ds := \frac{1}{2\pi i} \int_{\mathbb{R}} \varphi(q^* + i\tau) x^{-q^*} x^{-\tau} d\tau.$$ 

Lemma 2.4. If $f$ is in $L^1_{q}$ for every real number $q$ in $(a, b)$ then its Mellin transform is holomorphic in the strip $\{s \in \mathbb{C} \mid a + 1 < \text{Re } s < b + 1\}$.

Example 2.5.

(1) The function $f(x) = e^{-x}$ is in $L^1_{q}$ for all $q > -1$ so $Mf(s) = \Gamma(s)$ is holomorphic in the half-plane $\{s \in \mathbb{C} \mid \text{Re } s > 0\}$. 






The function $f(x) = \frac{1}{e^{\pi x} - 1}$ is in $L^1_q$ for all $q > 0$ so $Mf(s) = \Gamma(s)\zeta(s)$, where $\zeta(s) = \sum_{n=1}^{\infty} n^{-s}$ is the Riemann zeta function, is holomorphic in the half-plane $\{s \in \mathbb{C} \mid \Re s > 1\}$.

**Lemma 2.6.** If $f$ is in $L^1_q$ for every $q$ in $(a, b)$ then $Mf(\zeta + it)$ tends to 0 as $|t|$ tends to $+\infty$ uniformly with respect to $\zeta \in [a + 1 + \eta, b + 1 - \eta]$, for any $0 < \eta < \frac{b - a}{2}$.  

**Proposition.** For a function $f$ in $L^1_{q-\varepsilon} \cap L^2_{q+\varepsilon}$ we have:

$$\|f\|_{L^2_{q-\varepsilon}} = (2\pi)^{-1/2} \|Mf\|_{L^2_q}.$$

This identity may be directly derived from the Plancherel theorem for the Fourier transform of $g(y) = f(e^y)e^{\gamma y}$.

**Theorem (Plancherel transform).** According to the previous formula the Mellin transform can be extended, in a unique manner, to an isometry (up to the multiplicative constant $(2\pi)^{-1/2}$) of $L^2_q$ onto $L^2_{\tilde{q}}(q + i\mathbb{R})$, where $\tilde{q} = 2\pi^{1/2}$.

A consequence of the previous construction of the Mellin transform in $L^2$ is that for a function $f \in L^1_{q-\varepsilon} \cap L^2_{q+\varepsilon}$ (resp. $L^1_{q+\varepsilon}(q + i\mathbb{R}) \cap L^2_{q}(q + i\mathbb{R})$) the two definitions of the Mellin transform (resp. the inverse Mellin transform) coincide.

### 3. Solution of equation $\mathcal{L}(u) = f$ for self-similar kernels

We are now equipped to give a rigorous meaning to formula (14), which leads to solutions $u$ of equation (11).

The following assumption is equivalent to: $M^{-1}\left(\frac{Mf(s)}{1-k\mathcal{M}_0(s)}\right) \in L^1_q$. This assumption is not necessary in the $L^2_q$ case.

**Assumption 2.** $q \in \mathbb{R}$ and $f \in L^1_q$ are such that $Mf(s) \in L^1_q(q+i\mathbb{R})$ with $q^* = q + 1$ and there exists a function $v \in L^1_q$ such that $Mf(s)(1-k\mathcal{M}_0(s))^{-1} = Mv(s)$ for a.e. $s \in q + 1 + i\mathbb{R}$.

The link between the inverse Mellin transforms defined in $L^p_q$ for different values of $q$ is made via the zeros of the function $1 - k\mathcal{M}_0$, so we need an assumption on the location of the zeros of this function. The following assumption will be used for both the $L^1_q$ and the $L^2_q$ cases.

**Assumption 3.** For a given bandwidth $(a, b)$ with $a < 2 \leq b$, let $c \in (a, 2)$ such that the zeros of $1 - k\mathcal{M}_0$ in $\{s \in \mathbb{C} \mid a \leq \Re s \leq b\}$ are all included in the strip $\{s \in \mathbb{C} \mid c \leq \Re s \leq 2\}$ (such a real number $c$ exists as $\mathcal{M}_0$ is holomorphic and non-constant, cf (5)).

We assume that: there exists a set of horizontal segments

$$L_n = \{s \in \mathbb{C} \mid \Im s = c_n \text{ and } c \leq \Re s \leq 2\},$$

where $c_n$ tends to $\pm\infty$ when $n$ tends to $\pm\infty$, such that $|1 - k\mathcal{M}_0|$ is bounded from below by a positive constant on $\bigcup_{n \in \mathbb{Z}} L_n$.

#### 3.1. Holomorphic Mellin transform in $L^1$

All the proofs in this section can be found in appendix B and appendix C, and strongly rely on the results recalled in section 2.

The following proposition gives an explicit formula for $\mathcal{L}^{-1} : L^1_q \to L^1_q$ under some additional integrability assumptions.
Proposition 3.1. Let \( q > \text{Abs} \kappa_0 - 1 \) such that \( \zeta = q^+ = q + 1 \) satisfies assumption 1. For every \( f \) in \( L^1_q \) satisfying assumption 2, there exists a unique \( u \in L^1_q \) solution to equation (11). It is explicitly given for a.e. \( x > 0 \) by:

\[
u(x) := \frac{1}{2\pi i} \int_{q^+ - i\infty}^{q^+ + i\infty} \frac{Mf(s)}{1 - kM\kappa_0(s)} x^{-s} \, ds.
\] (17)

If \( f \) is in \( L^1_q \) for different values of \( q \), does formula (17) define the same solution \( u \), whatever the value of \( q \) is? And if not, how can we characterize their difference? Remember that in [14] (proposition A1), for the case \( \kappa_0 = \delta_{1/2} \), \( k = 2 \), there were exactly two distinct solutions, one holding for any \( q < 3 \) and one for any \( q > 3 \). The following proposition generalizes this result.

Proposition 3.2. Let \( a, b \) be such that \( \text{Abs} \kappa_0 < a \leq 2 < b \). We assume that the set \( P \) of the zeros of \( 1 - kM\kappa_0 \) satisfies assumption 3 and that assumption 1 is satisfied for any \( \zeta = q^+ \in (a, c) \cup (2, b) \). Let \( f \) satisfy assumption 2 for any \( q \in (a - 1, c - 1) \cup (1, b - 1) \).

For any \( q \in (a - 1, c - 1) \cup (1, b - 1) \), there exists a unique \( u_q \in L^1_q \) solution of equation (11), given by the explicit formula (17). The solution \( u_q = u_r(x) \) is independent of \( q \) in \( (a - 1, c - 1) \) (resp. \( u_q = u_r(x) \) is independent of \( q \in (1, b - 1) \)), and for almost every \( x > 0 \) we have:

\[
u_l(x) - u_r(x) = \frac{1}{2\pi i} \sum_{p \in P} a_p Mf(p) x^{-p} \quad \text{with} \quad a_p = \text{Res} \left( \frac{1}{1 - kM\kappa_0(s)}, s = p \right).
\]

Remark 3.3. Let \( g = (1 - kM\kappa_0(s))^{-1} \). From the relation \( M\kappa_0(s) = M\kappa_0(3) \) we deduce that the poles of \( g \) are pairwise conjugated. Also, if \( g \) can be represented as the Laurent series \( g(s) = \sum_{n \in \mathbb{Z}} l_n(s - p)^n \) near \( p \) then \( g \) can be represented as the Laurent series

Figure 1. Proofs of propositions 3.2, 3.10. The crosses are the zeros of \( 1 - kM\kappa_0 \).
Example 3.6. Consider $\kappa_0 = 1_{(0,1]}, k = 2$ and a function $f$ satisfying assumption 2 for $q$ in $(0, 2) \setminus \{1\}$. The functions $\mathcal{M} \kappa_0(x) = \frac{1}{2}$ and $\mathcal{M} f$ are holomorphic in the strip $\{ s \in \mathbb{C} \mid 1 < \text{Re} \, s < 3 \}$. The function $1 - 2 \mathcal{M} \kappa_0$ only vanishes at $s = 2$ and $\text{Re} \, s (\frac{1}{2}, s = 2) = 2$. For every $q^* \in (1, 3)$, $\mathcal{M} \kappa_0(q^* + i t)$ tends to $0 \neq \frac{1}{2}$ as $|t|$ tends to $+\infty$. Proposition 3.2 and remark 3.4 show that, for almost every $x > 0$:

$$u_i(x) - u_e(x) = \frac{1}{\pi} x^2 \int_0^{+\infty} t f(it) \, dt.$$

Example 3.7. Consider the general mitosis kernel $\kappa_0 = \delta_{i/k}$ for an integer $k \geq 2$ and a function $f$ in $L^2_0$ for all $q$ in $(0, 2)$ satisfying assumption 2 for $q$ in $(0, 2) \setminus \{1\}$. The functions $\mathcal{M} \kappa_0(x) = \frac{1}{k^{1-j}}$ and $\mathcal{M} f$ are holomorphic in the strip $\{ s \in \mathbb{C} \mid 1 < \text{Re} \, s < 3 \}$. The set of the zeros of $1 - k \mathcal{M} \kappa_0(x)$ is $2 + \frac{2\pi k}{\ln k} \mathbb{Z}$; it is easy to show that $\text{Re} \, s (\frac{1}{2}, s = 2 + \frac{2\pi k}{\ln k}) = \frac{1}{2}$ for every integer $n \in \mathbb{Z}$. For every $q$ in $(0, 2) \setminus \{1\}$, $|\mathcal{M} \kappa_0(q^* + i t)| = k^{-q} \neq \frac{1}{2}$. A simple computation leads to $|1 - k \mathcal{M} \kappa_0(2 + i y)| = 2 \sin \left( \frac{\pi k y}{2} \right)$; denoting $c_n = \frac{2\pi k}{\ln k}$ for $n \in \mathbb{Z}$ we find that $|1 - k \mathcal{M} \kappa_0|$ is bounded from below by 2 in $\cup_{n \in \mathbb{Z}} [c_n, 1]$. Proposition 3.2 and remark 3.4 show that, for almost every $x > 0$:

$$u_i(x) - u_e(x) = \frac{1}{2\pi \ln k} x^2 \sum_{n \in \mathbb{Z}} x^{-\frac{2\pi}{\ln k} n} \mathcal{M} f \left( 2 + \frac{2\pi i}{\ln k} n \right).$$

We refer the reader to appendix D for a study of the zeros of $1 - k \mathcal{M} \kappa_0$ if $\kappa_0 = \rho_1 \delta_{\rho_1} + \rho_2 \delta_{\rho_2}$. 

3.2. Isometric Mellin transform in $L^2$

We now turn to the $L^2$ theory which is easier to handle in the context of inverse problems and in the direct continuation of the previous studies [26]. All the proofs of this section can be found in appendix C.

In the $L^1$ theory, we saw that the $L^1_q$ norm for the function was linked to the $L^1(q^* + i \mathbb{R})$ space for its Mellin transform, with $q^* = q + 1$. In the $L^2_q$ theory, the $L^2$ space of the function is linked to the space $L^2_\nu (\mathbb{Q}^* + i \mathbb{R})$ for its Mellin transform, with $\mathbb{Q} = \frac{a_2}{a_1}$, as expressed by the Plancherel transform recalled in section 2. The following proposition gives an explicit bound for $\| L^{-1} \|_{L^2_q \to L^2_{\nu}}$. 

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\textbf{Proposition 3.8.} Let $q$ be a real number such that $\tilde{q} = \frac{q+1}{2} > \text{Abs}\kappa_0$ and assumption 1 is satisfied for $\xi = \tilde{q}$. For every $f$ in $L^2_q$ there exists a unique solution $u$ in $L^2_{\tilde{q}}$ to equation (11) and we have the estimate:

$$\|L^{-1}\|_{L^2_{\tilde{q}} \rightarrow L^2_q} \leq \sup_{\rho \in \mathbb{R}} \left| 1 - k\mathcal{M}\kappa_0 \left( \frac{q+1}{2} + i\rho \right) \right|^{-1}.$$ 

\textbf{Remark 3.9.} A weak form of this proposition is the following. We have for all real numbers $t$: $|\mathcal{M}\kappa_0 \left( \frac{q+1}{2} + i\rho \right) | \leq \mathcal{M}\kappa_0 \left( \frac{q+1}{2} \right)$. Thus if: $\mathcal{M}\kappa_0 \left( \frac{q+1}{2} \right) < \frac{1}{k}$, which is true if $\kappa_0((0, 1)) > 0$ and $q > 3$ for instance—then:

$$\sup_{\rho \in \mathbb{R}} \left| 1 - \mathcal{M}\kappa_0 \left( \frac{q+1}{2} + i\rho \right) \right|^{-1} \leq 1 - \mathcal{M}\kappa_0 \left( \frac{q+1}{2} \right) \right|^{-1}.$$ 

This gives explicit expressions for $\|L^{-1}\|_{L^2_{\tilde{q}} \rightarrow L^2_q}$ thanks to the table of example 2.2.

As for the $L^1$ theory, let us now turn to the case when $f$ is in $L^2_q$ for several $q$. The following proposition corresponds to proposition 3.2 in the $L^1$-case.

\textbf{Proposition 3.10.} Let $a$ and $b$ be such that $\text{Abs}\kappa_0 < a \leq 2 < b$. We assume that the set $P$ of the zeros of $1 - k\mathcal{M}\kappa_0$ satisfies assumption 3 and that assumption 1 is satisfied for any $\xi = \tilde{q}$ in $(a, c) \cup (2, b)$. Let $f \in L^2_q$ for any $q$ in $(2a - 1, 2b - 1)$.

For any $q$ in $(2a - 1, 2c - 1) \cup (3, 2b - 1)$, there exists a unique $u_q$ in $L^2_q$ solution of equation (11). For almost every $x > 0$, $u_q(x)$ does not depend on $q$ in $(2a - 1, 2c - 1)$ (resp. on $q$ in $(3, 2b - 1)$).

\textbf{Remark 3.11.} Propositions 3.1, 3.2, 3.8 and 3.10 can be easily adapted to closed equations, such as the case of the conservative equation, where $\kappa_0$ is still defined by (5), but the system keeps only one fragment at each division, so that the number of fragments $k$ does not appear in front of the non-local term (2) $(k = 1)$ and $\lambda = 0$:

$$\left\{ \begin{array}{l}
\frac{\partial}{\partial x}(g(x)N_c(x)) + B(x)N_c(x) = \int_0^{+\infty} \kappa(x, y)B(y)N_c(y) \, dy = 0, \quad x \geq 0, \\
g(0)N_c(0) = 0, \quad N_c(x > 0) > 0, \quad \int_0^{+\infty} N_c(x) \, dx = 1.
\end{array} \right. \quad (18)$$

In this case the linear operator $L_c$ is replaced by $L_c = \text{Id} - \mathcal{K}_c$, so that $1 - \mathcal{M}\kappa_0$ replaces $1 - k\mathcal{M}\kappa_0$ in assumptions 1, 2 and 3, and the ‘pivot vertical line’ is $1 + i\mathbb{R}$ instead of $2 + i\mathbb{R}$. The exceptional spaces are $L^1$ and $L^2_q$ (resp. instead of $L^1$ and $L^2_q$). See also remark E.5.

\section*{4. Numerical simulations}

\subsection*{4.1. Protocol}

We assume that $k \in \mathbb{N}^*$, $g$, $\kappa$ are known. Given a measure $N^\alpha$ of $\mathcal{N}$, we would like to compute approximations of the birth rate $B$. For $\alpha > 0$, we define an approximation $D_\alpha(N^\alpha)$ of $D(\mathcal{N})$ by formula (13).

Assuming that $N^\alpha$ is in $L^2_q$ for $q$ in $[q_1, q_2]$, $D_\alpha(N^\alpha)$ satisfies assumptions 1 and 2 thanks to the regularity properties of the mollifier sequence $(\rho_\alpha)$. We can invert the operator $L$ using the explicit formula (17) of proposition 3.1 along the line $\frac{q+1}{2} + i\mathbb{R}$ for the two values $q = q_1, q_2$, and compute approximations $H_{q_1, \alpha} = L_{q_1}^{-1}(D_\alpha(N^\alpha)) \in L^2_{q_1}$ and $H_{q_2, \alpha} = L_{q_2}^{-1}(D_\alpha(N^\alpha)) \in L^2_{q_2}$ of $H = BN$.
For a fixed $a > 0$ we then compute $H_\alpha(a)$ by formula (16).

Technological aspects. In order to compute integrals, we define the support of a function $f : \mathbb{R} \to \mathbb{C}$ as the set $f^{-1}([-f] \geq \eta])$ where $\eta$ is a small positive constant. The integral of a compactly supported function is computed with a closed Newton–Cotes formula (namely the Boole–Villarceau’s rule of order 5, cf [9]). A computation with the fast Fourier transform is much supported function is computed with a closed Newton–Cotes formula (namely the Boole–Villarceau’s rule of order 5, cf [9]). A computation with the fast Fourier transform is much faster (for $n$ points of integration only $O(n \ln n)$ operations are done instead of $O(n^2)$ with this method) but the results obtained are less accurate.

We denote by $F f(\xi) = \int_{\mathbb{R}} f(x) e^{-i\xi x} \, dx$, the Fourier transform of a function $f$. Thanks to the relation $F \rho \alpha(\xi) = F \rho(\alpha \xi)$, we compute the term $\rho \alpha * D(N^\epsilon)$ as:

$$(\rho \alpha * (gN^\epsilon))' + \lambda \rho \alpha * N^\epsilon = F^{-1}[F(\alpha \xi)(-i\xi F(gN^\epsilon)(\xi) + \lambda F(N^\epsilon)(\xi))].$$

The Mellin transform, and its inverse, are also linked to the Fourier transforms thanks to the relations:

$$M f(c + it) = F g(u) \quad \text{where} \quad g(u) = f(\exp(-u)) \exp(-cu),$$

$$M^{-1}_c f(x) = \frac{1}{2\pi} x^{-c} F h(\ln(x)) \quad \text{where} \quad h(z) = f(c + iz).$$

The terms appearing in proposition 3.1 are computed in this way.

Before applying our method on biological data, see section 4.3, we first tested it on simulations. We initially computed numerical approximations of the function $N$ using numerical schemes of [14, 15], and then added a random noise to $N$. Let us denote by $u^\epsilon$ the uniform random variable on the interval $\frac{1}{2}[-1, 1]$; we defined $N^\epsilon = (\max(N_k(1 + u^\epsilon), 0))_{k \in \{1, \ldots, n\}}$ where $N = (N_1, \ldots, N_n)$ is the numerical approximation of the density $N$. We then applied the previously described technique to get a numerical approximation of $H$.

4.2. Results on simulated data

Parameters. For all the following simulations we selected the growth rate $g(x) = x^{1/2}$, and the birth rate to reconstruct $B(x) = x^2$. The self-similar kernels tested satisfy $k = 2$ and are defined by $\kappa_0 = I_{[0,1]} \, dx$, $\delta_{1/2}$ or $\frac{1}{2\pi} \exp \left( -\frac{1}{2\sigma^2} \left( x - \frac{1}{2} \right)^2 \right) I_{[0,1]} \, dx$ with $\sigma = 0.1$. As proved theoretically in section 3.1, the functions $1 - 2M \kappa_0$, for $\kappa_0 = I_{[0,1]} \, dx$, $\delta_{1/2}$, only vanish

![Figure 2](image-url) - Reconstructions of $H$ and $B$ from exact numerical data $N$ for the truncated normal distribution and $\rho = \rho_1$. Dotted blue line: $H^1$, $B^1$; green dashed line: $H^2$, $B^2$; red line: $H^3$, $B^3$; dashed-dotted black line: $H = BN$, $B$. 


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Figure 3. Reconstructions of $H$ and $B$ from noisy numerical data $N'$ for $\kappa_0 = I_{[0,1]} dx$, $k = 2$, $(\varepsilon, \alpha) = (10\%, 32\%)$ and $\rho = \rho_2$.

Figure 4. $\ln \|H^3 - H^3\|_{L^2}$ as a function of $\ln \varepsilon$ for $\kappa_0 = \delta_{1/2}$, $k = 2$, and for $\rho = \rho_1$ (left), $\rho = \rho_2$ (right).

Simulations and figures for each kernel $\kappa_0$

(1) First no random noise is added to the numerically computed function $N$. Thus, we perform the Mellin inversion along the two vertical lines $0.5 + i\mathbb{R}$ and $3.5 + i\mathbb{R}$. This gives two numerical approximations of $H$, respectively $H^1$ in $L^2$ and $H^2$ in $L^2$. We also define $H^3 = H^1 I_{[0,2]} + H^2 I_{[2, +\infty)}$ which is in $L^q_{\kappa}$ for every $q$ in $[0, 6]$. We can then approximate the birth rate $B$ by defining $B' = H'/\max(N, \vartheta)$ where $\vartheta$ is a threshold, see [12]. See figure 2.
(2) Then a random noise is added to the numerically computed function $N$ to get a noisy datum $N^\varepsilon$ and we build numerical approximations—$H^\varepsilon_q \in L^2, H^\varepsilon_q \in L^2, H^\varepsilon_q \in L^2$ for every $q$ in $[0, 6]$—as previously. See figure 3.

The mean convergence rates for $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$, $i = 1, 2, 3$ and $\rho = \rho_1, \rho_2$ are the following; see also figure 4.

| $\rho$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ |
|-------|--------------------------------|--------------------------------|--------------------------------|
| $I_{[0, 1]} dx$ | $12.26 \varepsilon_{0.73}$ | $110 \varepsilon_{0.62}$ | $24 \varepsilon_{0.62}$ |
| $\delta_{1/2}$ | $5.76 \varepsilon_{0.61}$ | $90.15 \varepsilon_{0.68}$ | $23.21 \varepsilon_{0.68}$ |
| $G$ | $15.23 \varepsilon_{0.73}$ | $91.77 \varepsilon_{0.68}$ | $22.53 \varepsilon_{0.68}$ |
| $\rho_1$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ | $\|H^\varepsilon_q - H^\varepsilon_q\|_{L^2}$ |
| $I_{[0, 1]} dx$ | $5.22 \varepsilon_{0.76}$ | $7.74 \varepsilon_{0.72}$ | $5.77 \varepsilon_{0.75}$ |
| $\delta_{1/2}$ | $4.94 \varepsilon_{0.76}$ | $5.09 \varepsilon_{0.75}$ | $4.76 \varepsilon_{0.81}$ |
| $G$ | $4.85 \varepsilon_{0.76}$ | $5.36 \varepsilon_{0.73}$ | $3.84 \varepsilon_{0.76}$ |

All the illustrations are perfectly coherent with the theoretical results. The function $H^2 \in L^2$, which is a bad approximation of $H = BN$ near 0, whereas it is a better approximation of $H$ than $H^1$ near $+\infty$. The interpolated approximation $H^3$ is a good approximation of $H$ from both sides. The convergence rates of ($H^\varepsilon$) obtained numerically are better than the expected theoretical results (for $\rho_1$: $O(\varepsilon^{0.6})$ instead of $O(\varepsilon^{1/2})$; for $\rho_2$: $O(\varepsilon^{0.7})$ instead of $O(\varepsilon^{1/3})$). Note that the constants numerically obtained for $\rho_2$ are smaller than the ones for $\rho_1$.

4.3. Application to biological data

We analysed a dataset obtained through a microscopic time-lapse imaging of Escherichia coli cells—data published in [29]. Since the cells divide into two almost equal daughter cells, this corresponds to $k = 2$, and for the self-similar kernel $\lambda_0$, to an experimentally measured probability density $\kappa_0$ (see figure 5, right). The growth rate of the cells is $g(x) = cx$, with a value of $c$ that can be experimentally measured—we took here $c = 1$, which is always possible, up to a change in the time-scale. In this case the parameter $\lambda$ is equal to $c$ so here $\lambda = 1$. The experimental conditions are well-controlled and stable, so that the cells are in a steady state of growth. An image is taken every two minutes, and the length of each cell is measured through image analysis—since the cells are of cylindric shape with a roughly constant radius, we can identify the volume of a cell with its length, see [11, 27]. In the dataset we analysed a sample of 30,900 sizes, gathering all sizes at all times. As in [12], we model this sample as a $K$-sample of i.i.d. realizations of random variables of density $N(x)$, and we apply the mollifier directly to the dataset by defining

$$D_\rho(N^\varepsilon)(x) := \rho_\varepsilon \ast D \left( \frac{1}{K} \sum_{j=1}^{K} \delta_{x-x_j} \right),$$

see figure 5, left. The graph of $D_\rho(N^\varepsilon)$ looks like the curve of the Gaussian $N(\mu, \sigma)$ with $\mu = 0.49$ and $\sigma$ small. The mollifier is the probability density function of the normal distribution $N(0, 1)$.

We also are able to measure the daughters’ cell sizes just after division; we apply the same regularization technique to obtain the probability density $\kappa_0$. This fragmentation kernel is regularized with the optimal parameter for estimating normal densities.

To test our method, we first simulate the inverse problem with the regularized experimentally obtained steady state $D_\rho(N^\varepsilon)$ and the fragmentation kernels $\kappa_0 = \delta_{1/2}$ (see figure 6) and $\kappa_0$ obtained experimentally (see figure 7); we obtain the interpolated birth rate...
Figure 5. Regularized density law $D_\alpha(N^\epsilon)$ as a function of the sizes $x$ (left), with $\alpha = 5.6$, and regularized experimentally measured fragmentation kernel $\kappa_0$ (right), with $\alpha = 0.056$.

Figure 6. $H_i^{\epsilon,\alpha}$ and $B_i^{\epsilon,\alpha} = \frac{H_i^{\epsilon,\alpha}}{N^\epsilon}$ for $i = 1, 2, 3$ as functions of the sizes $x$ for the optimal value of $\alpha$ and for $\kappa_0 = \delta_{1/2}$, $\alpha = 0.81$.

Figure 7. $H_i^{\epsilon,\alpha}$ and $B_i^{\epsilon,\alpha} = \frac{H_i^{\epsilon,\alpha}}{N^\epsilon}$ for $i = 1, 2, 3$ as functions of the sizes $x$ for the optimal value of $\alpha$ and for $\kappa_0$ which is experimentally measured, $\alpha = 0.83$. 
Figure 8. Density law $D_\alpha(N^\varepsilon)$ and steady states $N(B_3^\varepsilon)$ for the optimal value of $\alpha$, for $\kappa_0 = \delta_{1/2}$ (left) and $\kappa_0$ experimentally measured (right). Discrepancies: 2.1% and 1.7%.

Second, to measure the accuracy of this method, we simulate the direct problem with $B_3^\varepsilon$ and the two fragmentation kernels $\kappa_0$ (see figure 8); we obtain the steady state $N(B_3^\varepsilon)$. To compare this steady state to $D_\alpha(N^\varepsilon)$ we define the discrepancy as:

$$\frac{\|D_\alpha(N^\varepsilon) - N(B_3^\varepsilon)\|_{L^2}}{\|D_\alpha(N^\varepsilon)\|_{L^2}}.$$

We choose the parameter $\alpha$, for the regularization of $N^\varepsilon$, which minimizes the discrepancy. For both kernels the optimal value for $\alpha$ is the same: $\alpha = 5.6$.

The discrepancies for both fragmentation kernels $\kappa_0$ are small, around 2%. The discrepancy is slightly smaller for the experimentally measured fragmentation kernel than for $\delta_{1/2}$.

5. Conclusion

We focussed on the problem of finding the birth rate $B$ of a size-structured population from measurements of the time-asymptotic profile $N^\varepsilon$, following [14, 15, 26]. For general fragmentation kernels we generalized the estimates obtained in [15] to every $L^p$ space for the birth rate $B$ and for less smooth data $N^\varepsilon$. Above all, for self-similar fragmentation kernels we provided a new method to find the birth rate $B$. The method is based on the Mellin transform, and is simple and efficient both theoretically and numerically. The method allowed us to understand the partial results for general fragmentation kernels and to improve them for self-similar kernels. The results obtained numerically confirmed the theoretical estimates. The method made it possible to compute the birth rate $B$ from a fragmentation kernel which is experimentally measured. The question of inverting the operator $L$ defined by (8) in $L^p(x^q \, dx)$ spaces for $q < 3$ remains open when the kernel $\kappa$ is not self-similar. Other research directions concerning adapting our method, based on the measurement of the asymptotic profile, to other types of structured population equations, possibly defined for multidimensional variables (cf [17, 24]). As concerns the numerical implementations, our method could probably be improved to run faster, which would be of key importance for higher dimension problems, using the fast Mellin transform for instance.
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Appendix A. The filtering method: proof of proposition 1.1

The filtering method is classical both in a deterministic setting and in a statistical setting, where it is generally referred to as the kernel estimation method. We give here a quick proof for its application in our setting, for the sake of completeness. This method gives a control of \( \| H^\varepsilon - H \|_{L_q^p} \) if the operator \( L : L_q^p \rightarrow L_q^p \) is invertible. We slightly generalize the statements of [15], where it was assumed that \( N^s \) is in \( L^2_{N^s} \), to a noisy measure \( N^s \in \mathbb{W}^{s}((1 + x^2) \, dx) \) spaces, which allows us to compare our results directly with those obtained in a statistical setting (heuristically comparable to \( s = -1/2 \), see [12]).

We first recall some convolution estimates, which are the basis of kernel regularization methods. The following lemma is a kind of generalization of lemma 2.1. of [15].

Lemma A.1. Let \( p \geq 1, q \geq 0, \alpha \in [-1, 1] \setminus \{0\}, \theta \in [0, 1] \) be real numbers and \( \rho \) a function supported in \([-1, 1]\). We define: \( \rho_\alpha(x) = \frac{1}{\alpha} \rho\left(\frac{x}{\alpha}\right) \) for \(|x| \leq \alpha \).

1. If the function \( f \) is in \( \mathbb{L}^p((1 + x^2) \, dx) \) and if \( \rho^{(n)} \) is in \( \mathbb{L}^1([-1, 1]) \) for an integer \( n \in \mathbb{N} \), we have:

\[
\left\| (\rho_\alpha)^{(n)} \ast f \right\|_{L_q^p} \leq C |\alpha|^{-n} \left\| f \right\|_{\mathbb{L}^p((1 + x^2) \, dx)} ,
\]

where \( C = \max(1, 2^{\frac{1}{p'}}) \left\| \rho^{(n)} \right\|_{\mathbb{L}^1([-1, 1])} \).

2. If the function \( f \) is in \( \mathbb{W}^{s}((1 + x^2) \, dx) \) and if \( \rho, \rho' \) are in \( \mathbb{L}^1([-1, 1]) \), we have:

\[
\left\| \rho_\alpha \ast f \right\|_{L_q^p} \leq C |\alpha|^{-\theta} \left\| f \right\|_{\mathbb{W}^{s}((1 + x^2) \, dx)} ,
\]

where \( C = \max(1, 2^{\frac{1}{p'}}) \max(\left\| \rho \right\|_{\mathbb{L}^1([-1, 1])}, \left\| \rho' \right\|_{\mathbb{L}^1([-1, 1])}) \).

3. If the function \( f \) is in \( \mathbb{W}^{s}((1 + x^2) \, dx) \) and if \( \rho', \rho'' \) are in \( \mathbb{L}^1([-1, 1]) \), we have:

\[
\left\| (\rho_\alpha)^' \ast f \right\|_{L_q^p} \leq C |\alpha|^{-1+\theta} \left\| f \right\|_{\mathbb{W}^{s}((1 + x^2) \, dx)} ,
\]

where \( C = \max(1, 2^{\frac{1}{p'}}) \max(\left\| \rho' \right\|_{\mathbb{L}^1([-1, 1])}, \left\| \rho'' \right\|_{\mathbb{L}^1([-1, 1])}) \).

4. If \( f \) is in \( \mathbb{L}^{1,p}((1 + x^2) \, dx) \) and if \( \rho \) is in \( \mathbb{L}^1([-1, 1]) \) with \( \left\| \rho \right\|_{\mathbb{L}^1([-1, 1])} = 1 \) we have:

\[
\left\| f - \rho_\alpha \ast f \right\|_{L_q^p} \leq C |\alpha| \left\| f \right\|_{\mathbb{W}^{s}((1 + x^2) \, dx)} ,
\]

where \( C = \max(1, 2^{\frac{1}{p'}}) \).

5. Let \( n \in \mathbb{N} \) an integer. If \( f \) is in \( \mathbb{W}^{n+1,p}((1 + x^2) \, dx) \) and if \( \rho \) satisfies \( \left\| \rho \right\|_{\mathbb{L}^1([-1, 1])} = 1 \) and \( \int_{-1}^1 \rho(z) \, dz = 0 \) for \( k = 1, \ldots, n \) we have:

\[
\left\| f - \rho_\alpha \ast f \right\|_{L_q^p} \leq C |\alpha|^{n+1} \left\| f \right\|_{\mathbb{W}^{n+1,p}((1 + x^2) \, dx)} ,
\]

where \( C = \max(1, 2^{\frac{1}{p'}})/(n!) \).
Proof of lemma A.1. We use the standard inequalities, valid for all real numbers $x, y$ and $\alpha \geq 0$:

$$\alpha \geq 1 \quad |x|^\alpha + |y|^\alpha \leq (|x| + |y|)^\alpha \leq 2^{\alpha - 1} (|x|^\alpha + |y|^\alpha),$$

$$\alpha \in [0, 1] \quad 2^{\alpha - 1} (|x|^\alpha + |y|^\alpha) \leq (|x| + |y|)^\alpha \leq |x|^\alpha + |y|^\alpha.$$

(1) We first prove this inequality in the case $n = 0$. Writing $|\rho| = |\rho|^\theta |\rho|^{1 - \theta}$ and by Hölder’s inequality, we have:

$$|\rho| * f(x) = \left| \int_{-1}^{1} \rho(z)f(x - \alpha z) \, dz \right| \leq \|\rho\|_{L^q([-1,1])}^{\frac{\theta}{p}} \left( \int_{-1}^{1} |\rho(z)|^p |f(x - \alpha z)|^p \, dz \right)^{\frac{1}{p}}.$$

Then:

$$\|\rho| * f\|_{L^p_x}^{p} \leq \|\rho\|_{L^q([-1,1])} \frac{1}{1 - \theta} \int_{-1}^{1} \int_{0}^{\infty} |f(x - \alpha z)|^p |\rho(z)|^q \, dx \, dz$$

$$= \|\rho\|_{L^q([-1,1])} \left( \int_{-1}^{1} |\rho(z)|^q \int_{0}^{\infty} |f(x - \alpha z)|^p |x - \alpha z|^q \, dx \, dz \right)^{\frac{1}{q}}$$

$$\leq \|\rho\|_{L^q([-1,1])} \left( \|f\|_{L^p_x}^{p} + |\alpha|^q \right) \left( \|\rho\|_{L^q([-1,1])} \right)^{-\frac{1}{q}}$$

$$\leq C \|\rho\|_{L^q([-1,1])} \|f\|_{L^p_x}^{p},$$

where $C = \max(1, 2^{\alpha - 1})$. The general case $n \in \mathbb{N}$ is a consequence of the previous inequality using the fact that $(\rho\alpha)^{n} = \alpha^n (\rho\alpha)^n$.

(2) We first prove the inequality for $\theta = 1$ and then use an interpolation argument to get it for $\theta$ in $[0, 1]$. Let $f = g + h$ be a function in $W^{-1,p}(1 + x^q) \, dx$ with $g, h$ in $L^p((1 + x^q) \, dx)$. Writing $\|\rho| * f\|_{L^p_x} \leq \|\rho| * g\|_{L^p_x} + \|\rho| * h\|_{L^p_x}$ applying the previous inequality for $n = 0$ and taking the infimum over $g, h$ such that $f = g + h$ we get the result. Let $\theta \in [0, 1]$ and let $T = \rho|*$. The first point of the lemma with $n = 0$ implies that $T$ is bounded from $L^p$ to $L^q((1 + x^q))$ with its norm less than $\max(1, 2^{\alpha - 1}) \|\rho\|_{L^q([-1,1])} \leq C$.

(3) We have just proved that $T$ is bounded from $L^p$ to $W^{-1,p}(1 + x^q) \, dx$ with its norm less than $C \frac{1}{q}$. Therefore by complex interpolation and by definition of $W^{-\theta,p}(1 + x^q) \, dx$ we get the estimate.

(4) As we have $f, f' \in L^p$ we can write: $f(x) - f(x - \alpha z) = \alpha \int_{0}^{1} f'(x - t\alpha z) \, dt$, so, by Hölder’s inequality or by Jensen’s inequality: $|f(x) - f(x - \alpha z)|^p \leq |\alpha|^p \int_{0}^{1} |f'(x - t\alpha z)|^p \, dt$. As $\|\rho\|_{L^q([-1,1])} = 1$ we have:

$$\|f - \rho| * f\|_{L^q_x} \leq \int_{0}^{\infty} \int_{-1}^{1} |f(x) - f(x - \alpha z)|^p \rho(z) \, dz \, dx$$

$$\leq \int_{0}^{\infty} \int_{-1}^{1} |f(x)|^p \rho(z) \, dz \, dx \leq \int_{-1}^{1} \int_{0}^{\infty} |f'(x - t\alpha z)|^p |\rho(z)| \, dz \, dx \, dt,$$

the rest of the proof is the same as that given in the first point of the lemma with $f'$ instead of $f$. 20
As we have $f, \ldots, f^{n+1} \in L^p$ we can write, for a.e. $x \in [0, +\infty)$:

$$f(x) - f(x - \alpha z) = \sum_{k=1}^{n+1} \frac{1}{k!} f^{(k)}(x - \alpha z)(\alpha z)^k + \frac{1}{n!} (\alpha z)^{n+1} \int_0^1 t^n f^{(n+1)}(x - t\alpha z) \, dt.$$  

Integrating this quantity multiplied by $\rho(z)$ we get:

$$f(x) - \rho_a * f(x) = \frac{\alpha^{n+1}}{n!} \int_0^1 z^{n+1} \rho(z) \int_0^1 t^n f^{(n+1)}(x - t\alpha z) \, dt \, dz.$$  

By the same manipulation as for the previous point we get:

$$\|f - \rho_a * f\|_{L^p} \leq \frac{\alpha^{n+1}}{(n!)^p} \int_0^1 \int_0^1 \int_0^{+\infty} |f^{(n+1)}(x - t\alpha z)|^p \rho(z) |x|^q \, dx \, dt \, dz,$$

which leads to the conclusion as previously. \hfill \square

**Proof of proposition 1.1.** Between $D(N)$ defined by (9) and $D_a(N)$ by (13), we define an intermediate $D_a$ by

$$D_a = \rho_a * D(N).$$

We use the decomposition: $\|D_a(N^\theta) - D(N)\|_{L^Q} \leq \|D_a(N^\theta) - D_a(N)\|_{L^Q} + \|D_a(N) - D(N)\|_{L^Q}.$

(1) For the term $\|D_a(N^\theta) - D_a(N)\|_{L^Q}$ it remains to estimate the two terms: $\|\rho_a * (gN' - gN)\|_{L^Q}$ and $|\lambda| \|\rho_a * (N' - N)\|_{L^Q}$ which can be done using respectively the third point of the lemma for $\theta = s$ and the second point for $\theta = s$.

(2) For the term $\|D_a(N) - D(N)\|_{L^Q}$ we write:

$$\|\rho_a * (gN') - (gN')\|_{L^Q} \leq \|\rho_a * (gN') - (gN')\|_{L^Q} + |\lambda| \|\rho_a * N - N\|_{L^Q},$$

and we use the fourth point or the fifth point of the lemma.

For $a, b > 0$, $\alpha, \varepsilon > 0$ we write: $\varepsilon \alpha^{-a} + \alpha^b = \alpha^{(b-a)/2} 2 \varepsilon^{1/2} (\varepsilon^{1/2} \alpha^{-a})^{1/2} + \varepsilon^{-1/2} \alpha^{(b+a)/2}$; the second term is greater than 2 with equality iff $\alpha = \varepsilon^{1/(b+a)}$. Thus we obtain the inequality:

$$\varepsilon \alpha^{-a} + \alpha^b \geq 2 \varepsilon^{b/(b+a)},$$

which is an equality iff $\alpha = \varepsilon^{1/(b+a)}$. The last points of the proposition are the consequences of this inequality with $(a, b) = (1 + s, 1)$ and $(a, b) = (1 + s, n + 1).$ \hfill \square

**Appendix B. Proofs of the lemmas of section 2**

**Proof of lemma 2.1.** This is a consequence of Lebesgue’s dominated convergence for holomorphic functions. For every real number $x > 0$: $s \mapsto x^{-s}$ is holomorphic in $C$ and for every $s \in C$ such that $\text{Re} \ s \geq \alpha$: $|x^{-s}| \leq x^{-1} \in L^1([0, 1], \mu).$ \hfill \square

**Proof of lemma 2.3.** The proofs are based on Lebesgue’s dominated convergence theorem and on Hölder’s inequality.

(1) If $a > b$, for all $x$ in $[0, 1]$: $x^a \leq x^b$. If $k = a(k_1 + (1 - a)k_2)$ with $\alpha$ in $(0, 1)$ the inequality $\mathcal{M}_0(k) \leq \mathcal{M}_0(k_1)\alpha^{a-b} \mathcal{M}_0(k_2)^{1-a}$ is a consequence of Hölder’s inequality for $p = \frac{1}{2} \in [1, +\infty]$.  

(2) For $x$ in $[0, 1]$, thus $\mu(a)$ because $\mu([1]) = 0$, $x^a$ tends to 0 as $a$ tends to $+\infty$ and is dominated by 1 everywhere.

(3) If $a > b$, for all $x$ in $(0, 1)$: $x^a < x^b$; and $0^0 = 0^0$, $1^a = 1^b$. From the fact that $\mu([0, 1]) > 0$ the result follows. \hfill \square
Proof of lemma 2.4. This is a consequence of Lebesgue’s dominated convergence theorem for holomorphic functions. For almost every real number \( x > 0 \), \( s \mapsto f(x) \zeta^{x-1} \) is holomorphic in \( \mathbb{C} \). Let \( \eta > 0 \). For all \( s \in \mathbb{C} \) such that \( a + 1 + \eta \leq \text{Re } s \leq b + 1 - \eta \) and for almost every \( x > 0 \): 
\[
|f(x)\zeta^{x-1}| \leq |f(x)| \zeta^{x} e^{\pm \eta} \mathbf{1}_{\{0,1\}}(x) + |f(x)| \zeta^{b-\eta} e^{\pm \eta} \mathbf{1}_{\{1,\infty\}}(x) \in L^1. 
\]
(In particular: \( \forall n \in \mathbb{N}^* \forall s \in \mathbb{C} \) such that \( a + 1 < \text{Re } s < b + 1 \), \( Mf^{(n)}(s) = f_0^{(n)}(\ln x) \zeta^{x-1} f(x) \text{ dx} \). \( \square \)

Proof of lemma 2.6. Let \( 0 < \eta < \frac{b-a}{4} \). For \( \varphi = \mathbf{1}_{[c,d]} \), with \( 0 < c \leq d \), \( M\varphi \) is holomorphic in \( \mathbb{C} \) using lemma 2.4 and can be explicitly computed \((M\varphi)(s) = \frac{1}{2} (c^2 - d^2)\) for \( s \neq 0 \) and \( M\varphi(0) = \ln \left( \frac{c}{d} \right) \); an easy computation leads to the desired property. By linearity of \( M \) it is also true for a linear combination of indicator functions.

Let \( \varepsilon > 0 \). There exists a linear combination of indicator functions \( \varphi \) such that \( \|f - \varphi\|_{L^1((0,1), x^{a+\eta} \text{ dx})} \leq \varepsilon \) and \( \|f - \varphi\|_{L^1((1,\infty), x^{b+\eta} \text{ dx})} \leq \varepsilon \). Therefore
\[
|M(f)\zeta^{i \theta}| \leq \int_0^1 |(f - \varphi)(x)| x^{a+\eta} \text{ dx} + \int_1^{+\infty} |(f - \varphi)(x)| x^{b-\eta} \text{ dx} + |M(\varphi)(\zeta + i \theta)|
\]
and then \( \lim_{|\zeta| \to +\infty} |M(f)\zeta^{i \theta}| \leq 2\varepsilon \) uniformly with respect to \( \zeta \) in \( [a + 1 + \eta, b + 1 - \eta] \). \( \square \)

Appendix C. Proofs related to subsections 3.1 and 3.2
Proofs 3.2 and 3.10 are illustrated by figure 1.

Proof of proposition 3.1. (Uniqueness) Let \( q^* = q + 1 \) and \( Z = 1 - kMk_0 \). The difference \( w \in L^1 \zeta \) of two solutions of the equation \( \mathcal{L}u = f \) satisfies \( \mathcal{L}w = 0 \) so \( \mathcal{L}(s) Mw(s) = 0 \) for all \( s \in q^* + i\mathbb{R} \). As \( \mathcal{L} \) does not vanish on \( q^* + i\mathbb{R} \), it implies that \( Mw = 0 \) on \( q^* + i\mathbb{R} \), thus \( w(x) = 0 \) for almost every \( x > 0 \) by the \( L^1 \) inversion theorem, cf theorem 2.2 section 2.

(Explicit formula) Let \( g(s) = Mf^{(1)}(s) \). Since \( Mf \in L^1(q^* + i\mathbb{R}) \), and since thanks to assumption 1, \( |Z^{-1}| \) is bounded (by \( 1/\delta \)) on \( q^* + i\mathbb{R} \), \( g \) is also in \( L^1(q^* + i\mathbb{R}) \). We can thus define a function \( u \) by the explicit formula (17).

(Existence) Thanks to assumption 2 there exists a function \( v \) in \( L^1 \zeta \) such that \( Mf(s)(Z(s))^{-1} = Mv(s) \) on \( q^* + i\mathbb{R} \) and then \( u = v \) a.e., by the \( L^1 \) inversion theorem, thus the previous explicit formula (17) defines an \( L^1 \zeta \) function. \( \square \)

Proof of proposition 3.2.
(a) Let \( g(s) = \frac{Mf(s)}{Z(s)} \), which is meromorphic in \( \{ s \in \mathbb{C} \mid a < \text{Re } s < b \} \) using proposition 2.4, and may have poles in \( P \). It results from proposition 3.1 that equation (11) \( \mathcal{L}(u) = f \) has a unique solution \( u_q \in L^1 \zeta \) for \( q \) in \( (a - 1, c - 1) \cup (1, b - 1) \) and that it is given for almost every \( x > 0 \) by \( u_q(x) = \frac{1}{2\pi i} \int_{q-i\infty}^{q+i\infty} g(s) x^{-s} \text{ ds} \).

(b) Let \( q_1, q_2 \in (1, b - 1) \). Denoting \( q^*_1 = q_1 + 1 \), we want to show that, for almost every \( x > 0 \):
\[
\int_{q^*_1 - i\infty}^{q^*_1 + i\infty} g(s) x^{-s} \text{ ds} = \int_{q_1 - i\infty}^{q_1 + i\infty} g(s) x^{-s} \text{ ds} + \int_{q_2 - i\infty}^{q_2 + i\infty} g(s) x^{-s} \text{ ds}.
\]

For \( n \in \mathbb{N} \) consider a rectangle \( \Gamma_n \) oriented in the direct sense with edges at \( q^*_1 \pm in \), \( q_1 \pm in \), and \( q_2 \pm in \). As the function \( Z \) does not vanish in the strip \( \{ s \in \mathbb{C} \mid q^*_1 \leq \text{Re } s \leq q_1 \} \) we have:
\[
0 = \int_{\Gamma_n} g(s) x^{-s} \text{ ds} = \int_{q_1 - in}^{q_1 + in} \cdots \text{ ds} + \int_{q^*_1 - in}^{q^*_1 + in} \cdots \text{ ds} + \int_{q_2 - in}^{q_2 + in} \cdots \text{ ds}.
\]
By lemma 2.6, as \( n \) tends to \(+\infty\) the two horizontal integrals \( \int_{q_1^*-ic_n}^{q_1^*+ic_n} g(s)x^{-t}\,ds \) and \( \int_{q_1^*-ic_n}^{q_1^*+ic_n} g(s)x^{-r}\,ds \) tend to zero (if \(|Mf(q^*)\| \leq \varepsilon\) for \(|t| \geq c_n, q^* \in [q_1^*, q_2^*]\), these integrals are smaller than \( \frac{1}{2} |q_2-q_1| |x|^{-q_2^*-q_1^*} \varepsilon \) where \( x^{a,b} := x^{a_1,1} + x^{b_1,1} \). By the dominated convergence theorem the two vertical integrals tend to the aimed integrals because \( \int_{q_1^*-ic_n}^{q_1^*+ic_n} |Mf(s)|\,ds \) and \( \int_{q_1^*-ic_n}^{q_1^*+ic_n} |Mf(s)|\,ds \) are finite. The same proof holds to show that \( u_\delta \) does not depend on \( q \) in \((a-1, c-1)\).

(c) To show that \( u_1 - u_2 \) is given by the above formula we use a path \( \Gamma_n \) encircling the set \( P_n \) set \( s_1 \in (a-1, c-1) \) and \( s_2 \in (1, b-1) \) and consider a rectangle oriented in the direct sense with edges at \( q_1^* \pm ic_n, q_2^* \pm ic_n \). The residue formula gives:

\[
\sum_{p \in \Gamma_n} \text{Res} \left( g(s)x^{-t}, s = p \right) = \int_{q_1^*-ic_n}^{q_1^*+ic_n} \cdots \int_{q_2^*-ic_n}^{q_2^*+ic_n} \int_{q_2^*-ic_n}^{q_2^*+ic_n} \cdots dx.
\]

The horizontal integrals tend to 0 as \( n \) goes to \(+\infty\) because of lemma 2.6 and because \(|s| \to 1 - kMk_0(t + ic_n)\) is uniformly bounded from below for \( t \) in \([q_1^*, q_2^*]\) and \( n \) in \( \mathbb{N} \) by assumption 3. The vertical integrals converge for the same reason as above.

\[ \square \]

**Proof of proposition 3.8.** (Uniqueness) Let \( \tilde{q} = \frac{q+1}{2} \). The difference \( w \in L^2_q \) of two solutions of equation (11) satisfies \( \mathcal{L}(w) = 0 \) so \( \mathcal{Z}(s)\mathcal{M}w(s) = 0 \) for all \( s \) in \( q + i\mathbb{R} \). As \( \mathcal{Z} \) does not vanish in \( q + i\mathbb{R} \), it implies that \( \mathcal{M}w = 0 \) in \( \tilde{q} + i\mathbb{R} \), thus \( w(x) = 0 \) for almost every \( x > 0 \) by theorem 2.2 of section 2.

(Existence) We define: \( \mathcal{M}u(s) = \mathcal{M}(f(s)\mathcal{Z}(s))^{-1} \) for all \( s \) in \( q + i\mathbb{R} \). As \( \mathcal{Z}^{-1} \) is bounded (by \( 1/\delta \)) on \( q + i\mathbb{R}, \mathcal{M}u \) is in \( L^2(\tilde{q} + i\mathbb{R}) \) and thus, using theorem 2.2 of section 2, \( u \) is in \( L^2_q \).

Moreover, the function \( u \) satisfies \( \mathcal{M}(f(s)) = \mathcal{M}(u(s))\mathcal{Z}(s) \) and is the solution of equation (11).

(Estimate) Using theorem 2.2 of section 2 we get the estimate:

\[
\|f\|_{L^2_q}^2 = \frac{1}{2\pi} \left\| \mathcal{M}(u(s)) \mathcal{Z}(s) \right\|_{L^2(\tilde{q} + i\mathbb{R})}^2 \leq \left\| \mathcal{Z}(s)^{-2} \right\|_{L^\infty(q + i\mathbb{R})} \frac{1}{2\pi} \left\| \mathcal{M}(u(s)) \right\|_{L^2(\tilde{q} + i\mathbb{R})}^2 \leq \|\mathcal{Z}\|_{L^\infty(q + i\mathbb{R})}^2 \|u\|_{L^2_q}^2.
\]

\[ \square \]

**Proof of proposition 3.10.** Let \( q \in (2a - 1, 2b - 1) \) and \( \tilde{q} = \frac{q+1}{2} \in (a, b) \). Proposition 3.8 implies that equation (11) with \( f \in L^2_q \) admits a unique solution \( u_\delta \in L^2_q \) for \( q \) in \((2a - 1, 2c - 1) \cup (3, 2b - 1) \) and that it is defined by \( u_\delta = \mathcal{M}^{-1}(\frac{\mathcal{M}_f}{\mathcal{M}_g}) \), with \( \mathcal{Z} = 1 - kMk_0 \), and where the inverse is taken on the line \( q + i\mathbb{R} \).

Let \( T > 0 \). We define:

\[
f_T = f \mathbb{1}_{[\frac{1}{T}, \frac{T}{2}]} \text{ on } [0, +\infty), \quad \varphi_T(s) = \exp \left( \frac{s^2}{T} \right), \quad g_T(s) = \frac{\mathcal{M}f_T(s)}{\mathcal{Z}(s)} \varphi_T(s) \text{ on } \mathbb{C}.
\]

The function \( f_T \) is in \( L^1_{\tilde{q}} \) for every \( q \in \mathbb{R} \) thanks to Hölder’s inequality. As a consequence the function \( \mathcal{M}f_T \) is holomorphic on \( \mathbb{C} \), is in \( L^\infty(q + i\mathbb{R}) \) for every \( q \in \mathbb{R} \) and the function \( \mathcal{M}f_T \varphi_T \) is in \( L^1(q + i\mathbb{R}) \) for every \( q \in \mathbb{R} \) because

\[
|\varphi_T(q + iy)| \leq \exp \left( \frac{q^2 - y^2}{T} \right) \in L^\infty;
\]

the function \( g_T \) is meromorphic in \( \{ s \in \mathbb{C} \mid a^* < \text{Re} \, s < b^* \} \) and may have poles in \( P \).

Let \( q_1 \) and \( q_2 \) be in \((3, 2b - 1) \). We want to show \( u_{q_1}(x) = u_{q_2}(x) \) for almost every \( x > 0 \).

As in point (b) of the proof of proposition 3.2, taking paths along the lines \( \tilde{q}_1 + i\mathbb{R} \) and \( \tilde{q}_2 + i\mathbb{R} \),
Appendix D. Study of the zeros of $1 - k \mathcal{M}((\rho_1 \delta_{\sigma_1} + \rho_2 \delta_{\sigma_2})$ along vertical lines

Consider $\kappa_0 = \rho_1 \delta_{\sigma_1} + \rho_2 \delta_{\sigma_2}$, with $0 < \sigma_1 < \sigma_2 < 1$ and $\sigma_1 \leq \frac{1}{2} \leq \sigma_2$. The Mellin transform of $\kappa_0$ is $\mathcal{M}\kappa_0(s) = \rho_1 \sigma_1^{s-1} + \rho_2 \sigma_2^{s-1}$ for $s \in \mathbb{C}$. The assumptions $\mathcal{M}\kappa_0(1) = \rho_1 + \rho_2 = 1$ and $\mathcal{M}\kappa_0(2) = \rho_1 \sigma_1 + \rho_2 \sigma_2 = \frac{3}{2}$ are satisfied iff $(\rho_1, \rho_2) = \left(\frac{\sigma_2-\sigma_1}{\sigma_2-\sigma_0}, \frac{1}{\sigma_2-\sigma_0}\right)$. The measure $\kappa_0$ is non-negative because $\sigma_1 \leq \frac{1}{2} \leq \sigma_2$.

For a given real number $x$, we want to describe the sets $\Gamma_x = \{x + iy \mid y \in \mathbb{R}, \rho_1 \sigma_1^{s-1} \sigma_1^y + \rho_2 \sigma_2^{s-1} \sigma_2^y = \frac{1}{2}\}$, which are the zeros of $1 - k \mathcal{M}(\rho_1 \delta_{\sigma_1} + \rho_2 \delta_{\sigma_2})$ on the line $x + i\mathbb{R}$. Each set $\Gamma_x$ is discrete because $\mathcal{M}\kappa_0$ is a non-constant holomorphic function. As $\kappa_0((0, 1)) > 0$, the function $t \in \mathbb{R} \mapsto 1 - k \mathcal{M}\kappa_0(t) \in \mathbb{R}$ is strictly increasing, cf lemma 2.3, so $2 \in \Gamma_2$ (because $\mathcal{M}\kappa_0(2) = \frac{3}{2}$), $\Gamma_x$ is the empty set for $x > 2$ and for every $x < 2$, $\Gamma_x \cap \mathbb{R}$ is the empty set. Describing the set $\Gamma_x$ is based on the following geometrical fact.

Let $A$, $B$ and $C$ be non-negative real numbers. Finding a pair of real numbers $(\theta, \theta')$ satisfying $A e^{\theta} + B e^{\theta'} - C = 0$ is equivalent to building a positively oriented triangle, possibly flat, with sides of lengths $A, B, C$ and such that the oriented angles between the

\[\begin{align*}
\text{Figure D1.} & \quad \text{The triangle condition.}
\end{align*}\]
sides of lengths $C$ and $A$ and the sides of lengths $C$ and $B$ are respectively equal to $\theta$ and $\theta'$. If $\Theta_0 = (\theta_0, \theta'_0)$ is a solution of the previous equation then all the solutions are $\Theta_0 + 2\pi \mathbb{Z} \times 2\pi \mathbb{Z}$ $\cup$ $-\Theta_0 + 2\pi \mathbb{Z} \times 2\pi \mathbb{Z}$ (in the complex plane the problem is to find the intersection of the circle centred at 0 of radius $A$ with the circle centred at $C$ of radius $B$).

Let $x \in \mathbb{R}$ and set $A = \rho_1 \sigma_1^{-1}, B = \rho_2 \sigma_2^{-1}, C = \frac{1}{5}$. The equation $A\sigma_1^0 + B\sigma_2^0 = C$ has a solution iff one can find $y \in \mathbb{R}$ such that a triangle as above can be constructed with $\theta = y \ln \sigma_1$ mod $2\pi$ and $\theta' = y \ln \sigma_2$ mod $2\pi$ (triangle condition). If there exists a solution $y = y_0$ then $y$ is another solution iff $(y - y_0)(\ln \sigma_1, \ln \sigma_2) \in (2\pi \mathbb{Z})^2$ iff $y \in y_0 + 2\pi \mathbb{Z}$ where $G$ is the additive subgroup of $\mathbb{R}$, $G = \alpha \mathbb{Z} \cap \beta \mathbb{Z}$ with $\alpha = \frac{1}{5\sigma_1}$ and $\beta = \frac{1}{5\sigma_2}$. The structure of such a group is discussed below. If there exists $y_0$ such that the triangle condition is true then $\Gamma_y = x + iH$ with $H = y_0 + 2\pi G \cup -y_0 + 2\pi G$. Otherwise $\Gamma_y = \emptyset$.

By the triangle condition, in order to have a solution we must have $|A - B| \leq \frac{1}{5} \leq A + B$, where $A = \rho_1 \sigma_1^{-1}, B = \rho_2 \sigma_2^{-1}$. The right-hand side of this inequality is not satisfied for $x > 2$, and we find again that $\Gamma_y$ is the empty set in this case. In the case $x = 2$ the triangle condition is satisfied with $y = 0$ so $\Gamma_2 = 2 + 2\pi iG$ where $G$ is the same as above. In the case $x < 2$, the triangle condition may be satisfied or not (for instance if $|A - B| > \frac{1}{5}$).

To finish describing the sets $\Gamma_y$, we have to study the subgroup of $\mathbb{R}$, $G = \alpha \mathbb{Z} \cap \beta \mathbb{Z}$, where $\alpha$ and $\beta$ are two given real numbers. The subgroups of $\mathbb{R}$ are either cyclic, that is of the form $c\mathbb{Z}$ for some real number $c \geq 0$, or are dense in $\mathbb{R}$. Therefore, as $G$ is discrete, it is cyclic. It is easy to prove that there exists $c > 0$ such that $c\mathbb{Z} \cap c\mathbb{Z} = c\mathbb{Z}$ if $\alpha \neq 0$ and $\beta \neq 0$ and $\frac{\alpha}{\beta} \in \mathbb{Q}$.

We then have two cases. The first one is when $\frac{\alpha}{\beta}$ is a rational number (for instance if $\sigma_2 = 1 - \sigma_1$ and $\sigma_1 = \frac{1}{\sqrt{2}}, r = \frac{\sqrt{3}}{3}, r' = \frac{\sqrt{3} + 1}{3}$); then $G = c\mathbb{Z} \cap c\mathbb{Z}$ is cyclic and not reduced to $\{0\}$. The set $\Gamma_y$ is $x + iH$ or is the empty set whether the triangle condition is satisfied or not. The second (and generic) case is when $\frac{\alpha}{\beta}$ is not a rational number (for instance if $\sigma_2 = 1 - \sigma_1$ with $\sigma_1$ transcendental); then $G$ is reduced to $\{0\}$ and consequently $\Gamma_y = \{2\}$ and for $x < 2$, $\Gamma_y = x \pm i y_0$ or $\Gamma_y = \emptyset$ whether the triangle solution is satisfied or not.

Appendix E. Inverting the operator $\mathcal{L}$ as a Neumann series

It is possible to build directly $\mathcal{L}^{-1} : L^p \rightarrow L^p$ for any $p \geq 1$ and for large enough weights $q$ ($q > 2p - 1$) thanks to a Neumann series. This result generalizes proposition 2.1 of [15], which corresponds to the case $p = 2$ and $q > 3$, and was proved by the use of Young’s inequalities and the Lax–Milgram theorem.

The following constants in $\mathbb{R} \cup \{\pm \infty\}$ are defined in [15].

Definition E.1. Let $\kappa$ be a general fragmentation kernel. We set, for a real number $u$:

$$C(u) = \sup_{x \in (0, \infty)} \int_x^{\infty} \kappa(x, y) \left(\frac{y}{x}\right)^u \, dy, \quad D(u) = \sup_{y \in (0, \infty)} \int_0^y \kappa(x, y) \left(\frac{y}{x}\right)^u \, dx.$$ 

Remark E.1. In the case of a self-similar kernel $\kappa(x, y) = \frac{1}{y} \kappa_0 \left(\frac{x}{y}\right)$ the constants of definition E.1 are simply written:

$$C(u) = M \kappa_0(u), \quad D(u) = M \kappa_0(u + 1).$$

Proposition E.2. Let $p, r$ be two real numbers with $1 < p < \infty$. The linear operator $\mathcal{K} : L^p \rightarrow L^p$ defined by (7) has a norm smaller than

$$C_r := C \left(\frac{q + 1}{p} - r\right)^{1/p} D \left(\frac{q + 1}{p} + r(p - 1) - 1\right)^{1/p}.$$
where $C(u), D(u)$ are introduced in definition E.1 and where $p'$ is the conjugate exponent of $p$ (that is: $\frac{1}{p} + \frac{1}{p'} = 1$).

If there exists a real number $r$ such that $C_r < \frac{1}{k}$, the linear operator $\mathcal{L} : \mathbb{L}_q^p \rightarrow \mathbb{L}_q^p$ defined by (8) is invertible and we have the estimate:

$$
\|\mathcal{L}^{-1}\|_{\mathbb{L}_q^p \rightarrow \mathbb{L}_q^p} \leq (1 - kC_r)^{-1}.
$$

**Proof.** Let $u \in \mathbb{R}$. Writing $\kappa = \kappa^{1/p}(\frac{x}{y})^{u/p'}\kappa^{1/p}(\frac{y}{x})^{-u/p'}$ for $x, y > 0$ and using Hölder’s inequality we have:

$$
\int_x^{+\infty} \kappa(x, y) |H(y)| \, dy 
\leq \left( \int_x^{+\infty} \kappa(x, y) \left( \frac{x}{y} \right)^w \, dy \right)^{1/p'} \left( \int_x^{+\infty} \kappa(x, y) \left( \frac{x}{y} \right)^{-uw/p'} |H(y)|^p \, dy \right)^{1/p}.
$$

Then using Fubini’s theorem for non-negative functions:

$$
\|K\|_{\mathbb{L}_q^p} \leq C(u)^{\frac{2}{p'}} \int_0^{+\infty} x^p \int_x^{+\infty} \kappa(x, y) \left( \frac{x}{y} \right)^{-uw/p'} |H(y)|^p \, dy \, dx
= C(u)^{\frac{2}{p'}} \int_{y>0} |H(y)|^p y^q \int_0^y \kappa(x, y) \left( \frac{x}{y} \right)^{-uw/p'} \, dx \, dy
\leq C(u)^{\frac{2}{p'}} D(-u(p - 1) + q) \|H\|_{\mathbb{L}_q^p}.
$$

This calculation leads to the conclusion by taking $r = \frac{q + 1}{p'} - u$. \hfill \Box

**Remark E.3.** If we have $f$ in $\mathbb{L}_q^p$ we can define the affine map $T : \mathbb{L}_q^p \rightarrow \mathbb{L}_q^p$ by $T(u) = f + kK(u)$. By the previous calculations $T$ is Lipschitz continuous with its constant less than $kC_r$ and is therefore a contraction if this constant is strictly less than one. Let us call $u$ the unique fixed point of $T$, which is the solution of $L u = f$. Defining the sequence $u_1 = T(u_0)$ for some $u_0$ in $\mathbb{L}_q^p$, we have that $(u_n)$ tends to $u$ with geometric convergence rate: $\|u_{n+1} - u\|_{\mathbb{L}_q^p} \leq \frac{(kC_r)^n}{1-kC_r} \|u_1 - u_0\|_{\mathbb{L}_q^p}$.

This remark was used to implement a code to compute $L^{-1} f$ for a given $f$ in $\mathbb{L}_q^p$. The inverse obtained in this way is in $\mathbb{L}_q^p$ for large enough values of $q$.

**Corollary E.4.** Let $p$ be a real number with: $1 < p < \infty$ and let $\kappa(x, y) = \frac{1}{x} \kappa_0(\frac{x}{y})$ be a self-similar kernel and $k \in \mathbb{N}^*$ defined by (5). The linear operator $\mathcal{L} : \mathbb{L}_q^p \rightarrow \mathbb{L}_q^p$ defined by (8) is invertible if $M_{\kappa_0}(\hat{q}) < \frac{1}{k}$ where $\hat{q} = \frac{q + 1}{p}$ with the estimate: $\|L^{-1}\|_{\mathbb{L}_q^p \rightarrow \mathbb{L}_q^p} \leq (1 - kM_{\kappa_0}(\hat{q}))^{-1}$. Moreover, if $\kappa_0((0, 1)) > 0$ the condition $M_{\kappa_0}(\hat{q}) < \frac{1}{k}$ is satisfied iff $q > 2p - 1$.

**Proof.** The first part results from proposition E.2 for $r = 0$ and from remark E.1. The second part is a consequence of lemma 2.3: $M_{\kappa_0}(\hat{q}) < \frac{1}{k} = M_{\kappa_0}(2)$ iff $\hat{q} > 2$ \hfill \Box

**Remark E.5.** For the conservative equation introduced in remark 3.11 with a self-similar kernel $\kappa(x, y) = \frac{1}{x} \kappa_0(\frac{x}{y})$ corollary E.4 becomes: the linear operator $\mathcal{L} = \text{Id} - K : \mathbb{L}_q^p \rightarrow \mathbb{L}_q^p$ defined by (7) is invertible if $M_{\kappa_0}(\hat{q}) < 1$ where $\hat{q} = \frac{q + 1}{p}$. If $\kappa_0((0, 1)) > 0$ the condition $M_{\kappa_0}(\hat{q}) < 1$ is satisfied iff $q > p - 1$.  

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For a self-similar kernel $\kappa_0$ such that $\kappa_0((0, 1)) > 0$ the condition of this corollary is true for any fixed $p$ in $(1, +\infty)$ and for large enough values of $q$, unlike the conditions obtained in propositions 3.1 and 3.8 which are satisfied only for $p = 1$ and 2 but for any weight $q$ except a countable set of values of $q$.

References

[1] Baccelli F, McDonald D R and Reynier J 2002 A mean-field model for multiple TCP connections through a buffer implementing red Perform. Eval. 49 77–97
[2] Balagué D, Cañizo J and Gabriel P 2013 Fine asymptotics of profiles and relaxation to equilibrium for growth-fragmentation equations with variable drift rates Kinetic Relat. Models 6 219–43
[3] Banks H T, Sutton K L, Thompson W C, Bocharov G, Roos E D, Schenkild T and Meyerhans A 2011 Estimation of cell proliferation dynamics using CFSE data Bull. Math. Biol. 73 116–50
[4] Bell G Land Anderson E C 1967 Cell growth and division: I. a mathematical model with applications to cell volume distributions in mammalian suspension cultures Biophys. J. 7 329–51
[5] Cáceres M J, Cañizo J A and Mischler S 2011 Rate of convergence to the remarkable state for fragmentation and growth-fragmentation equations J. Math. Pures Appl. 96 334–62
[6] Calvez V, Lenuzza N, Oels D, Deslys J-P, Laurent P, Mouthon F and Perthame B 2009 Size distribution dependence of prion aggregates infectivity Math. Biosci. 1 88–99
[7] Chiorino G, Metz J A J, Tomasoni D and Ubezio P 2001 Desynchronization rate in cell populations: mathematical modeling and experimental data J. Theor. Biol. 208 185–99
[8] Cloez B 2012 Limit theorems for some branching measure-valued processes arXiv:1106.0660v2
[9] Crouzeix M and Mignot A L 1984 Analyse Numérique des Équations Différentielles (Paris: Masson)
[10] Doumic M and Gabriel P 2009 Eigenelements of a general aggregation-fragmentation model Math. Models Methods Appl. Sci. 20 757
[11] Doumic M, Hoffmann M, Krell N and Robert L 2013 Statistical estimation of a growth-fragmentation model observed on a genealogical tree ArXiv:1210.3240
[12] Doumic M, Hoffmann M, Reynaud P and Rivoirard V 2012 Nonparametric estimation of the division rate of a size-structured population SIAM J. Numer. Anal. 50 925–50
[13] Doumic M, Maia P and Zubelli J P 2010 On the calibration of a size-structured population model from experimental data Acta Biotheoretica 58 405–13
[14] Doumic M, Perthame B and Zubelli J P 2009 Numerical solution of an inverse problem in size-structured population dynamics Inverse Problems 25 045008
[15] Doumic M and Tine L M 2013 Estimating the division rate for the growth-fragmentation equation J. Math. Biol. 67 69–103
[16] Engel H W, Hanke M and Neubauer A 1996 Regularization of Inverse Problems (Mathematics and its Applications vol 375) (Berlin: Springer)
[17] Delaunay F, Feillet C, Robert N, Billy F and Clairambault J 2013 Age-structured cell population model to study the influence of growth factors on cell cycle dynamics Math. Biosci. Eng. 10 1–17
[18] Groh A, Krebs J and Wagner M 2011 Efficient solution of an inverse problem in cell population dynamics Inverse Problems 27 065009
[19] Laurençot P and Perthame B 2009 Exponential decay for the growth-fragmentation/cell-division equation Commun. Math. Sci. 7 503–10
[20] Metz J A J and Diekmann O (ed) 1983 The Dynamics of Physiologically Structured Populations (Lecture Notes in Biomathematics vol 68) (Berlin: Springer) (Papers from the colloquium held in Amsterdam)
[21] Michel P, Mischler S and Perthame B 2005 General relative entropy inequality: an illustration on growth models J. Math. Pures Appl. 84 1235–60
[22] Pakdaman K, Perthame B and Salort D 2010 Dynamics of a structured neuron population Nonlinearity 23 55–75
[23] Pakdaman K, Perthame B and Salort D 2012 Adaptation and fatigue model for neuron networks and large time asymptotics in a nonlinear fragmentation equation submitted
[24] Pakdaman K, Perthame B and Salort D 2013 Relaxation and self-sustained oscillations in the time elapsed neuron network model SIAM J. Appl. Math. 23 1260–79
[25] Perthame B and Ryzhik L 2005 Exponential decay for the fragmentation or cell-division equation J. Diff. Eqs 210 155–77
[26] Perthame B and Zubelli J P 2007 On the inverse problem for a size-structured population model
Inverse Problems 23 1037–52
[27] Robert L, Hoffmann M, Krell N, Aymerich S, Robert J and Doumic M 2013 Division control in
escherichia coli is based on a size-sensing rather than timing mechanism submitted
[28] Rudin W 1962 Fourier Analysis on Groups (New York: Interscience)
[29] Stewart E J, Madden R, Paul G and François T 2005 Aging and death in an organism that reproduces
by morphologically symmetric division PLoS Biol. 3 e45
[30] Titchmarsh E C 1937 Introduction to the Theory of Fourier Integrals (Oxford: Clarendon)
[31] Villani C 2013 Cours de deuxième année donné à l’école normale supérieure de Lyon
http://cedricvillani.org/wp-content/uploads/2013/03/ana2.pdf 2003
[32] Wahba G 1977 Practical approximate solutions to linear operator equations when the data are noisy
SIAM J. Numer. Anal. 14 651–67