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Mathematical modeling of lymphocytes selection in the germinal center

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Abstract Lymphocyte selection is a fundamental process of adaptive immunity. In order to produce B-lymphocytes with a target antigenic profile, mutation selection and division occur in the germinal center, a specific part of lymph nodes. We introduce in this article a simplified mathematical model of this phenomenon, taking into account the main mechanisms. This model is written as a non-linear, non-local, inhomogeneous second order partial differential equation, for which we develop a mathematical analysis. We assess, mathematically and numerically, in the case of piecewise-constant coefficients, the performance of the biological function by evaluating the duration of this production process as a function of several parameters such as the mutation rate or the selection profile, in various asymptotic regimes.

Keywords Division-mutation-selection \cdot Germinal center \cdot Somatic hypermutation \cdot Affinity maturation \cdot Adaptive immunity \cdot Parabolic partial differential equation \cdot Population dynamics

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1 Introduction

Understanding the immune system is a key challenge in current biological, medical and pharmaceutical research, leading to revolutionary biomedical applications such as vaccination, immunotherapy or specific antibody production.

Adaptive immunity is responsible for the evolution of the antibody repertoire, a learning process necessary to identify and fight new foreign pathogens. This adaptation relies on a Darwinian process of Division-Mutation-Selection (DMS) occurring in the germinal centers (Kelsoe 1996), where an explosion of the mutation rate associated with B-cells division, called somatic hypermutation (Teng and Papavasiliou 2007), is observed, hence providing a unique example of an evolutionary process occurring within living organisms.

Although the general qualitative description of this process is well-established in the literature (Wabl et al. 1999; Neuberger et al. 1998; Diaz and Casali 2002), the quantitative assessment of this DMS process has remained largely unexplored experimentally,

in particular due to the difficulty to gather precise phylogenetic data of the B-cell repertoire during the various phases (Dunn-Walters et al. 2002). Recently, several biological studies (Gitlin et al. 2014; Victora 2014; Tarlinton 2014) have provided new experimental insights about the microscopic features of B-cell dynamics in the germinal center. Moreover, several key questions remain highly debated, such as the recycling of selected B-cells (Kepler and Perelson 1993; Meyer-Hermann et al. 2001) or the neutrality of the mutation process. Over the last few years, this system has been studied using relatively detailed computational models and numerical simulations Swerdlin et al. (2008), Meyer-Hermann (2002), providing a method to investigate several hypothesis and phenomena such as B-cell migration in the lymph node or the impact of recycling.

The aim of this article is to introduce a simplified macroscopic mathematical model of this process, in order to establish rigorous mathematical foundations and to investigate theoretically the impact of a few key parameters, such as the mutation rate or the selection profile, on the performance of the B-cell production, characterized by the duration of the process or the final quality of the repertoire.

Mathematical modeling of population dynamics can be approached either from a microscopic agent-based point of view, considering the behavior of many individuals, leading to a stochastic system in high dimension, or from a macroscopic point of view, where global quantities such as the number of individuals in a given state are considered, leading to partial differential equations (PDE) or integro-differential equations. In this article, we focus on the latter approach, and introduce a PDE model describing a population of B-cells subject to division, mutation and selection. In our framework, the division and mutation features give rise to classical linear diffusion terms, whereas selection introduces an inhomogeneous term. When a sufficient amount of B-cell with desired properties has been selected, the overall process shall terminate, which is modeled through a feedback term. As a result, the proposed model is a non-linear, non-local and inhomogeneous elliptic PDE.

Several mathematical models of evolutionary dynamics with mutation and selection have been previously studied, especially in population dynamics (Perthame 2006; Bürger 2000; Nowak 2006), ranging from the early works of Fisher (1999) to advanced mathematical models of adaptive dynamics. The model we introduce in this article has the particularity to combine spatial inhomogeities (in the space of traits) with a non-linear global feedback, which give rises to a specific PDE, for which we establish general properties, as well as precise estimates describing the impact of relevant parameters. In particular, our aim is to understand how the interplay between the mutation rate and the selection function influence the characteristic time-scale of the B-cell production process.

We first start with existence and uniqueness results for solutions of our new model. This gives a precise meaning to the solution, its regularity under the more general hypotheses. From the mathematical point of view, the main difficulty is due to the stiff, non-linear and non-local source term. The non-local feature is expressed in time and in the trait space.

In a second step, we compute, for a birth rate that is piecewise constant, the production time, *i.e*, the time t_{ϱ_0} for which ϱ_0 , a threshold selected population, is reached. This is performed with respect to two parameters of the model: ε the width of the selection window, and μ the mutation rate. The results of this study lead to several conclusions:

- For a general traits-square-integrable initial population of *B* cells,
 - If ε is small enough, and μ is greater than a threshold value, then t_{ϱ_0} behaves as $|\log \varepsilon|$, and is independent on μ ,
 - For a fixed ε , if μ is small, we observe that t_{ϱ_0} tends towards a constant that depends on the initial birth rate, ε and initial conditions.
 - For a fixed ε , when the mutation rate becomes large, t_{ϱ_0} tends to another finite value, still depending on the same parameters.
- For an initial datum which is a Dirac mass supported at a given trait z outside of the selection window:
 - If the domain is unbounded, we show that when μ tends to either 0 or ∞ , t_{ϱ_0} blows up,
 - On the other hand, when the domain is finite, t_{ϱ_0} stays bounded for μ growing large.

Intuitively these results can be interpreted by the paradigm: the larger the size of the repertoire, the more likely there should be an optimal mutation rate in terms of efficiency.

The distinction between different initial conditions comes from a debate in the literature, opposing supporters of mono-clonal germinal centers at the beginning of the process (Perelson and Weisbuch 1997) to authors trying to prove and measure oligo-clonal initial populations (Faro and Or-Guil 2013).

The article is organized as follows. In Sect. 2, we define mathematically our model and discuss its motivations, assumptions and limitations. In Sect. 3, we derive general results concerning existence and uniqueness of solutions, as well as their quantitative properties. To gain further understanding into the dynamical behavior of the system with respect to the data, we study in Sect. 4 the asymptotic behavior of t_{ϱ_0} when ε becomes small, whereas in 5 we consider for a fixed ε , the asymptotic regimes when μ is either large or small. This section focusses as well on different types of initial *B*cell population. Throughout the paper, theoretical results are illustrated with numerical simulations.

2 Mathematical model

We consider the time evolution of a population of lymphocytes during the Division-Mutation-Selection process within the germinal center.

2.1 Biological background

In this first section, we provide an elementary summary of the relevant biological background to describe our modeling approach. For the interested reader, we refer to Abbas et al. (2012) for a more detailed account.

The immune response to an external pathogen (virus, bacteria, etc.,) involves many different types of cells and employs various strategies to eliminate the pathogenic sources.

One of the most important way to fight pathogens relies on the bonding between antigens and antibodies, which triggers an efficient immune response, recruiting many other agents such as macrophages or T-lymphocytes. Antibodies are macromolecular compounds, made of peptidic chains, and whose purpose is to bond with antigens, which are complementary molecules presented at the surface of pathogens. The antibody-antigens (A–A) bonding can be thought as putting a key in a lock, and is characterized by the concept of affinity, which quantifies the likelihood that this bonding occurs for a specific A–A pair.

The production of antibodies in the immune response is ensured by the B-type lymphocyte, which is an immune cell able to produce a single specific antibody. Therefore, it is essential for B-cells to be able to learn how to produce high affinity antibodies with respect to a given antigen. Studying this learning process is precisely the purpose of the present article.

Once a new antigen, called the target, is detected by the immune system, it is captured by the follicular dendritic cells, and brought to the lymph nodes. A simplified vision of the process of affinity maturation in the germinal center, a specific part of the lymph node, can be summarized as follows:

- An initial population of a few immature B-cells enters in the germinal center. The affinity between the initial antibodies carried by these B-cells is in general relatively low, although it is not known whether this initial choice of B-cells is generic, or is already somehow adapted to the target.
- 2. During the first three days, B-cells divide and the population increases, until it migrates to another part of the germinal center, called the light zone. Notice that this population of B-cells is only able to produce a few types of different antibodies, typically less than 10.
- 3. In the light zone, the B-cells are now subject to a full division-mutation-selection process, which is the focus of our model.
 - Division- Mutation: At each division, a single B-cell produces two cells, one of which has undergone a significant amount of mutations in the part of its DNA responsible for the production of the antibody peptide sequence. This process is called somatic hypermutation, since the mutation rate is now increased to extremely high levels, several orders of magnitude higher than in normal cell divisions.
 - Selection: follicular dendritic cells are in the light zone and present at their surface the target antigen. B-cells are also presenting their antibody at their surface and wander in a seemingly random manner in the light zone. Then, the affinity between the presented antibody and the target antigen determines the probability that a bonding occurs and lasts long enough. In that case, the B-cell is selected and receives a signal which enables the cell to escape from the germinal center. Otherwise, low affinity B-cells, which were not able to receive such a signal, die.
- 4. After selection, a B-cell might have several fates: either it transforms into a plasma cell, which is able to produce and release antibodies in order to fight the pathogen, or into a memory cell whose aim is to remember the antigen (therefore being able to produce quickly high affinity antibodies, in case the pathogen returns later). A

third possible fate which is also discussed in the literature (Meyer-Hermann et al. 2001) is the possibility for the selected B-cell to come back inside the germinal center and to be subject a second time to the DMS process.

5. At some point, this affinity maturation terminates, and the precise biological mechanisms responsible for the determination of the stopping time remain unclear. However, it seems reasonable to consider that the process would stop as soon as a sufficient quantity of selected B-cell has escaped the germinal center.

2.2 Assumptions of the model

Of course, the above description is only a simplified and partial overview of a process which is in reality more complex and involves many different cell types, in particular T-cells. However, we think that it is neither possible nor wishful to take into account all the details of real biological processes into a mathematical model. Therefore, after running numerical simulations of various models (agent-based, stochastic models, PDE models with many variables), and motivated by the idea of introducing a simple, yet non-trivial, macroscopic mathematical model of the evolution of a population of B-cells during DMS phase, we have identified and selected what we consider to be the key parts of this complex process.

Space of traits: We consider that each B-cell is characterized by a *trait* corresponding to a specific antibody sequence. Instead of a discrete space of traits, composed with strings of amino-acids (as for instance in Balelli et al. 2016a, b), we view here the trait as an abstract property of the antibody and we therefore consider the space of traits to be made with real numbers, for instance the interval [0, 1].

Affinity and selection: The target antigen is also characterized by a trait x_0 in the same space, and if x denotes the trait associated with a B-cell, then we consider that the A-A affinity is given by a function $s(x) = F(x_0, x)$: the higher is s(x), the higher is the likelihood that the B-cell with trait x binds to the target. Since the affinity with the target summarizes the necessary information about a given B-cell to decide its fate, it seems reasonable to consider real trait x, directly translated into an affinity through the function *F*. However, this model does not address the difficult question of understanding how a small change in the DNA of the B-cell will result into a change in the peptide chain of the antibody and finally into a modification of the affinity.

Mutations: In our modeling approach, we consider that a mutation will change slightly the trait x to x + dx in a diffusion manner, and that the affinity will also change slightly through the function F. This model of mutation does not take into account the possibility for a small mutation of the DNA to produce a large change in affinity. One way to overcome this difficulty would be to consider a non-local mutation kernel instead of a diffusion, but it seems relatively uneasy to us to make precise and justified assumptions on such a kernel.

Termination: We assume that the termination of the affinity maturation process is regulated by a measure of the number of selected B-cells. More precisely, we assume that the division rate is a decreasing function of the number of selected B-cells. There-

fore, if this number reaches a certain value, the birth rate becomes lower than the death rate and the population inside the germinal center should extinct.

2.3 Mathematical model

We are now able to define precisely the mathematical model we will consider in this article. First, we introduce the following notations:

- 1. Parameters:
 - Cell division rate function $Q : \mathbb{R} \to \mathbb{R}_+$ is monotonic decreasing and $Q_1 \le Q(z) \le Q_0$
 - Cell death rate d > 0
 - Affinity-dependent selection function s(.), peaked around x_0 , the target.
 - Mutation rate μ > 0, which may either be a constant or a bounded function of the trait μ : ℝ → ℝ^{*}₊
- 2. Variables:
 - n(t, x) is the quantity of lymphocytes with a trait x
 - $-\rho(t)$ is the quantity of selected lymphocytes at time t and is given by:

$$\varrho(t) := \int_0^t \int_{\mathbb{R}} s(x) n(\tilde{t}, x) dx d\tilde{t}$$

The domain in the trait-space is denoted Ω and could be practically seen as a distance to a specific target trait. We are now able to formulate the main evolution equation, for x in the trait space $\Omega := (0, 1)$ and $t \ge 0$. It is an initial boundary value problem reading : find the function n solving

$$\begin{cases} \partial_t n(t,x) = (\mathcal{Q}(\varrho(t)) - d - s(x))n(t,x) + \partial_x(\mu\partial_x n(t,x)), & (t,x) \in \mathcal{O}_T := (0,T) \times \Omega, \\ \mu\partial_{\mathbf{n}} n(t,\cdot) = 0, & (t,x) \in \Sigma_T := (0,T) \times \partial\Omega, \\ n(0,x) = n_I(x), & \{0\} \times \Omega, \end{cases}$$
(1)

where the second line is the homogeneous Neumann boundary condition $(\partial_{\mathbf{n}} n := \partial_x n \cdot \mathbf{n})$, and the third one is the setting of initial data at time t = 0. Notice that in the case where μ is a constant scalar, then the main equation simply reads:

$$\partial_t n(t, x) = (Q(\varrho(t)) - d - s(x))n(t, x) + \mu \partial_{xx} n(t, x).$$

3 General results

In this section, we establish general results of existence and uniqueness, as well as spectral decomposition, concerning system 1 under various assumptions on the coefficients.

3.1 Existence and uniqueness

Here, we present general existence and uniqueness results concerning system (1) (Theorem 1 for a Lipschitz continuous Q and Theorem 2 for a piecewise constant Q).

We consider the Banach space

$$X = C([0, T]; L^{2}(\Omega)), \quad ||m||_{X} := \sup_{0 \le t \le T} ||m(t)||_{L^{2}(\Omega)},$$

for some T chosen later. Moreover we set

$$V(\mathcal{O}_T) := L^{\infty}((0,T); L^2(\Omega)) \cap L^2((0,T) \times H^1(\Omega))$$

and we define the form

$$I(t_1, n, \eta) := \int_{\Omega} n(t_1, x) \eta(t_1, x) dx - \int_0^{t_1} \int_{\Omega} n(t, x) \partial_t \eta(t, x) dx dt + \int_0^{t_1} \int_{\Omega} \mu \partial_x n(t, x) \cdot \partial_x \eta(t, x) + (s(x) + d - Q(\varrho)) n(t, x) \eta(t, x) dx.$$

Definition 3.1 We call a weak solution of problem (1) any solution $n \in V(\mathcal{O}_T)$ s.t.

$$I(T, n, \eta) = 0, \tag{2}$$

for every function $\eta \in H^1(\mathcal{O}_T)$ s.t. $\eta \equiv 0$ when t = 0. We say moreover that the solution is consistent with the initial condition if

$$I(T, n, \eta) = \int_{\Omega} n_I(x)\eta(0, x)dx, \quad \forall \eta \in H^1(\mathcal{O}_T).$$

the latter equation will be denoted as the variational formulation associated to the problem (1).

Hypotheses 3.1 Concerning the data, we assume

- (i) The initial data $n_I(x)$ belongs to $L^2(\Omega)$ and is non-negative.
- (ii) The function Q is a globally Lipschitz with respect to ρ , *i.e.* $Q \in W^{1,\infty}(\mathbb{R})$.
- (iii) The selection function s is a bounded non-negative function of x.

(iv) The mutation rate μ is a bounded positive definite function of x.

(v) The death rate d is a non-negative constant.

Theorem 1 Under hypotheses 3.1, there exists a unique positive weak solution $n \in V(\mathcal{O}_T)$ for any positive time T.

Proof We suppose in a first step that the constant d is strictly positive. Then we prove the existence using the Banach fixed point theorem. We define A a closed subset of X

$$A = \{m \in X, m > 0, \|m\|_X \le C_{\Phi}\},\$$

where C_{ϕ} is defined so that $||n_I||^2_{L^2(\Omega)} + ||Q||^2_{L^{\infty}}TC^2_{\phi}/(4d) < C^2_{\phi}$. We denote by $c_1 := ||Q||^2_{L^{\infty}}/(4d)$. For each $m \in A$, let *n* be the weak solution associated to the problem:

$$\begin{aligned} \partial_t n(t, x) &- \partial_x (\mu \partial_x n(t, x)) + (d + s(x))n(t, x) = \mathcal{Q}(\varrho(t))m(t, x), & (t, x) \in \mathcal{O}_T \\ \varrho(t) &= \int_0^t \int_{\Omega} s(x)m(\tilde{t}, x)dxd\tilde{t}, & t \in (0, T), \\ \partial_{\mathbf{n}} n(t, x) &= 0 & (t, x) \in \Sigma_T, \\ n(t = 0, x) &= n_I(x) > 0, & t = 0, x \in \Omega. \end{aligned}$$
(3)

This system defines the operator $\Phi : m \mapsto n$. We prove that it admits a unique point in *A*.

- Existence:

The solution $n \in V(\mathcal{O}_T)$ exists uniquely by standard parabolic theory (Theorem 5.1 p. 170 chap. III Ladyženskaja et al. 1968) for any given $m \in L^{\infty}(\mathcal{O}_T)$. The L^{∞} bound follows the same way by Theorem 7.1 p. 181 (Ladyženskaja et al. 1968). Moreover by Theorem 4.2 p. 160 we know that

$$\int_0^{T-h} \|n(t+h,\cdot) - n(t,\cdot)\|_{L^2(\Omega)}^2 dt = o(h).$$

- Non-negativeness:

We suppose that *m* is a positive function in $V(\mathcal{O}_T) \cap L^{\infty}(\mathcal{O}_T)$. We follow results from p. 183 Ladyženskaja et al. (1968) and we choose $n_h^- := \min(n_h, 0)$ as a test function in (2) where n_h is the Steklov approximation of *n* i.e.

$$n_h(t,x) = \frac{1}{h} \int_t^{t+h} n(\tau,x) \, d\tau, \quad \forall (t,x) \in (0, T-h) \times \Omega.$$

Applying such a test function is possible since n_h^- actually does belong to $H^1(\mathcal{O}_T)$. Passing to the limit with respect to the small parameter *h* then gives the identity

$$\frac{1}{2} \left[\int_{\Omega} (n^{-})^{2}(\tau, x) dx \right]_{\tau=0}^{\tau=t_{1}} + \int_{0}^{t_{1}} \int_{\Omega} \left\{ \mu \partial_{x} n \partial_{x} n^{-} + (s(x) + d) n n^{-} \right\} dx dt$$
$$= \int_{0}^{t_{1}} \int_{\Omega} Qmn^{-} dx dt.$$

But the support of n^- is the set where $n \le 0$ thus one has due to the positivity of m and Q that

$$\frac{1}{2} \left[\int_{\Omega} (n^{-})^2(\tau, x) dx \right]_{\tau=0}^{\tau=t_1} \le 0,$$

which gives that

$$\int_{\Omega} (n^{-})^2 (t_1, x) dx \le 0, \quad \forall t_1 \le T.$$

Since this integral is also non-negative, we deduce that the negative part of n vanishes, which proves the claimed non-negativeness property.

– The map Φ is an endomorphism in *X*:

By the same technique as above, we test by n_h and pass to the limit with respect to h (for a more detailed explanation see p. 141-142 Ladyženskaja et al. 1968) in the weak formulation, which writes:

$$\frac{1}{2} \left[\int_{\Omega} n^2(\tau, x) dx \right]_{\tau=0}^{\tau=t_1} + \int_{\mathcal{O}_{t_1}} \{ \mu |\partial_x n|^2 + (d+s(x))n^2 \} dx dt$$
$$= \int_{\mathcal{O}_{t_1}} \mathcal{Q}(\varrho(t)) m n \, dx dt,$$

and applying Cauchy-Schwarz and the Young inequalities on the right hand side gives

$$\frac{1}{2} \left[\int_{\Omega} n^2(\tau, x) dx \right]_{\tau=0}^{\tau=t_1} \le \frac{\|Q\|_{L^{\infty}} T}{4d} \|m\|_X^2 \le \frac{\|Q\|_{L^{\infty}} T}{4d} C_{\Phi}^2, \quad \forall t_1 < T.$$

Using the hypothesis on C_{Φ} we then deduce that in turn

$$\forall t \in [0, T], ||n(t, \cdot)||_{L^2} \leq C_{\Phi},$$

and thus $n \in A$.

– The map Φ is a contraction:

We denote $n_i = \Phi(m_i)$ for $i \in \{1, 2\}$ where m_i are two given functions in A. Then we denote $\tilde{n} := n_1 - n_2$ and $\tilde{m} := m_1 - m_2$, and we have

$$\begin{split} &\frac{1}{2} \left[\int_{\Omega} \tilde{n}^{2}(\tau, x) dx \right]_{\tau=0}^{\tau=t_{1}} + \int_{\mathcal{O}_{t_{1}}} \{ \mu | \partial_{x} \tilde{n} |^{2} + (d+s(x)) \tilde{n}^{2} \} dx dt \\ &= \int_{\mathcal{O}_{t_{1}}} (\mathcal{Q}(\varrho_{1}(t)) m_{1} - \mathcal{Q}(\varrho_{2}) m_{2}) \tilde{n} \, dx dt, \\ &\leq \frac{1}{4d} \int_{\mathcal{O}_{t_{1}}} \left[\mathcal{Q}(\varrho_{1}) m_{1} - \mathcal{Q}(\varrho_{2}) m_{2} \right]^{2} dx dt + d \int_{\mathcal{O}_{t_{1}}} \tilde{n}^{2} dx dt. \end{split}$$

We establish a bound for the first term in the latter right hand side

$$\frac{1}{4d} \int_{\mathcal{O}_{t_1}} \left[\mathcal{Q}(\varrho_1) m_1 - \mathcal{Q}(\varrho_2) m_2 \right]^2 dx dt \le \frac{t_1 \|\mathcal{Q}\|_{L^{\infty}}^2}{2d} \|\tilde{m}\|_X^2 + \frac{\|m_2\|_X^2}{2d} \int_0^{t_1} |\tilde{\mathcal{Q}}|^2 dt.$$

The last term above can then be estimated using that

$$\int_0^{t_1} |\tilde{Q}|^2 dt = \int_0^{t_1} |Q(\varrho_1) - Q(\varrho_2)|^2 dt \le \|Q\|_{\text{Lip}}^2 \|s\|_{L^2(\Omega)}^2 \|\tilde{m}\|_X^2 \frac{t_1^3}{3}$$

which gives

$$\|\tilde{n}(t_1, \cdot)\|_{L^2(\Omega)}^2 \le t_1(c_1 + t_1^2 C_{\Phi}^2 c_2) \|\tilde{m}\|_X^2$$

where

$$c_2 = \frac{\|Q\|_{L^{\infty}}^2}{2d} \quad c_3 = \frac{\|Q\|_{\text{Lip}}^2 \|s\|_{L^2(\Omega)}^2}{6d}$$

This finally provides

$$\|\tilde{n}\|_{X} \le \sqrt{T(c_{2} + C_{\Phi}^{2}T^{2}c_{3})} \|\tilde{m}\|_{X}$$

Choose T small enough so that the contraction holds. The local existence and uniqueness follow from the Banach-Picard theorem.

- Global existence

Denote $S_k = \sum_{i=1}^{k} T_i$. By induction, we assume that existence and uniqueness of (1) hold until the time S_k with the corresponding bound C_k . Now consider the time interval $[S_k, S_{k+1}]$, the new problem is the variational formulation corresponding to the system written in a strong form:

$$\begin{aligned} \partial_{t}\check{n}_{in}(t,x) &= \left[\mathcal{Q}(\varrho_{out}(t)) - d - s(x) \right] \check{n}_{in}(t,x) + \partial_{x}(\mu \partial_{x}\check{n}_{in})(t,x), & (t,x) \in (S_{k}, S_{k+1}) \times \Omega, \\ \varrho_{out}(t) &= \varrho_{out}(S_{k}) + \int_{[S_{k},t] \times \Omega} s(x)\check{n}(\tilde{t},x) dx d\tilde{t}, & t \in (S_{k}, S_{k+1}), \\ \check{n}_{in}(t = S_{k},x) &= n(S_{k},x), & (t,x) \in (\{S_{k}\} \times \Omega, \\ \partial_{n}\check{n}_{in}(t,0) &= 0. & (t,x) \in ((S_{k}, S_{k+1}) \times \partial\Omega. \end{aligned}$$

Iterating the same argument as in the first and second steps, existence and uniqueness hold on this new time interval if there exist $(T_k, C_k)_{k \in \mathbb{N}}$ s.t.

$$\begin{cases} \|n(S_k, \cdot)\|_{L^2}^2 + c_1 T_{k+1} C_{k+1}^2 < C_{k+1}^2, \\ \left(c_2 + C_{k+1}^2 T_{k+1}^2 c_3\right) T_{k+1} < 1, \end{cases} \quad \forall k \in \mathbb{N}. \end{cases}$$

The first condition ensures that the map Φ is an endomorphism while the second one insures that it is a contraction. Because $||n(S_k, \cdot)||_{L^2} \leq C_k$, the first inequality holds if we suppose that

$$C_k^2 + c_1 T_{k+1} C_{k+1}^2 < C_{k+1}^2.$$

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We choose (T_k) and (C_k) as

$$T_k := \frac{1}{2k c_1}, \quad C_k := \alpha k,$$

where α is s.t.

$$\alpha < \frac{c_1}{c_3}(4c_1 - c_2) = \frac{c_1}{2c_3d} \|Q\|_{L^{\infty}}^2$$

the series $S_k = \sum_{i=1}^k T_i$ diverges. Thus global uniqueness and existence hold.

For a fixed positive constant *d* we proved the theorem: there exists a unique weak solution $n_d \in V(\mathcal{O}_T)$, for any time *T*, solving

$$I_d(t_1, n_d, \eta) = 0, \quad \forall \eta \in H^1(\mathcal{O}_T) \quad \text{s.t. } \eta \equiv 0 \text{ for } t = 0.$$

Testing again with an appropriate averaged test function and passing to the limit gives:

$$\frac{1}{2} \left[\int_{\Omega} n_d^2(\tau, x) dx \right]_{\tau=0}^{\tau=t_1} + \int_{\mathcal{O}_{t_1}} \{ \mu | \partial_x n_d |^2 + s(x) n_d^2 \} dx dt \le \int_{\mathcal{O}_{t_1}} \mathcal{Q}(\varrho_d(t)) n_d^2 dx dt,$$

which provides by standard techniques

$$\|n_d\|_{V(\mathcal{O}_{t_1})} \leq C e^{\|Q\|_{L^{\infty}t_1}} \|n_I\|_{L^2(\Omega)}.$$

This estimate is uniform with respect to d. By weak convergence, one passes to the limit when $d \rightarrow 0$. This proves the theorem in this specific case.

Hereafter we weaken the Lipschitz hypothesis made on Q and we define

Hypotheses 3.2 We suppose that Hypotheses 3.1 (i),(iv) hold, moreover we suppose that

- (ii)' Q is a smooth function on $\mathbb{R} \setminus \{\varrho_0\}$, it admits two possibly different limits in the neighborhood of $\varrho_0 Q_{\pm} := \lim_{\varrho \to \varrho_0^{\pm}} Q(\varrho)$.
- (iii)' *s* is a positive definite function : $\inf_{x \in \Omega} s(x) > 0$.

Theorem 2 Under Hypotheses 3.2, one has the same conclusions as in Theorem 1, except for uniqueness that holds until a time t_0^- defined below.

Proof We define Q_{δ} a regularized non-linear source term

$$Q_{\delta}(\varrho) := \begin{cases} Q(\varrho) & \text{if } \varrho \in \mathbb{R} \setminus]\varrho_0, \varrho_0 + \delta[, \\ Q_- + \frac{Q(\varrho_0 + \delta) - Q_-}{\delta}(\varrho - \varrho_0) & \text{otherwise.} \end{cases}$$

As Q_{δ} is now Lipschitz continuous, one applies Theorem 1, denoting n_{δ} the corresponding unique solution, it belongs to $V(\mathcal{O}_T)$ for any time T, uniformly with respect

to δ . Moreover, one has also that $n_{\delta} \in C(0, T; L^2(\Omega))$ uniformly with respect to δ . Thus, one might extract a subsequence denoted again n_{δ} that converges weakly in $L^2(0, T; H^1(\Omega))$ and weak star in $L^{\infty}(0, T; L^2(\Omega))$. For every fixed $t \ge 0$, $n_{\delta}(t, x)$ tends to n(t, x) in $L^2(\Omega)$ weak. This shows that as $\delta \to 0$,

$$\begin{split} &\int_{\Omega} n_{\delta}(t_{1}, x)\eta(t_{1}, x)dx \to \int_{\Omega} n(t_{1}, x)\eta(t_{1}, x)dx, \\ &\int_{\mathcal{O}_{t_{1}}} n_{\delta}(t, x)\partial_{t}\eta(t, x)dxdt \to \int_{\mathcal{O}_{t_{1}}} n(t, x)\partial_{t}\eta(t, x)dxdt, \\ &\int_{\mathcal{O}_{t_{1}}} \mu\partial_{x}n_{\delta}(t, x)\partial_{x}\eta(t, x)dxdt \to \int_{\mathcal{O}_{t_{1}}} \mu\partial_{x}n(t, x)\partial_{x}\eta(t, x)dxdt, \\ &\int_{\mathcal{O}_{t_{1}}} (s(x) + d)n_{\delta}(t, x)\eta(t, x)dxdt \to \int_{\mathcal{O}_{t_{1}}} (s(x) + d)n(t, x)\eta(t, x)dxdt. \end{split}$$

It remains to show that

$$\int_{\mathcal{O}_{t_1}} \mathcal{Q}_{\delta}(\varrho_{\delta}) n_{\delta}(t, x) \eta(t, x) dx dt \to \int_{\mathcal{O}_{t_1}} \mathcal{Q}(\varrho) n(t, x) \eta(t, x) dx dt.$$

- ρ_{δ} converges strongly to ρ in C(0, T). Indeed, because $s \in L^{\infty}(\Omega)$,

$$\rho_{\delta}(t) = \int_0^t \int_{\Omega} s(x) n_{\delta}(t, x) dx dt \to \int_0^t \int_{\Omega} s(x) n(t, x) dx dt =: \rho(t).$$

and because s and n are non-negative functions ρ is a non-decreasing function. - ρ_{δ} is a strictly increasing

In the weak formulation $I(t, n_{\delta}, \eta) = 0$, we set η as a time dependent function independent on *x* that solves

$$\begin{cases} \partial_t \eta = (\|s\|_{L^{\infty}(\Omega)} + d - Q(\varrho_{\delta}(t)))\eta, & t \in (0, t_1) \\ \eta(0) = 1, & t = 0 \end{cases}$$

 η is then explicit and reads:

$$\eta = \exp\left(\int_0^t (\|s\|_{L^{\infty}(\Omega)} + d - Q(\varrho_{\delta}(\tilde{t})))d\tilde{t}\right).$$

We denote hereafter \overline{n} the average of *n* with respect to the trait variable :

$$\overline{n} := \int_{\Omega} n(t, x) dx.$$

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Finally one has that:

$$\overline{n}_{\delta}(t) \exp\left(\int_{0}^{t} (\|s\|_{L^{\infty}(\Omega)} + d - Q(\varrho_{\delta}(\tilde{t})))d\tilde{t}\right)$$
$$= \overline{n}_{\delta}(0) + \int_{\mathcal{O}_{t}} (\|s\|_{L^{\infty}} - s(x))n_{\delta}(\tilde{t}, x)\eta(\tilde{t})dx d\tilde{t}$$

the latter term being non-negative, and because the initial condition does not depend on δ , one concludes:

$$\begin{aligned} \overline{n}_{\delta}(t) &\geq \overline{n}(0) \exp\left(\int_{0}^{t} (\mathcal{Q}(\varrho_{\delta}(\tilde{t})) - (\|s\|_{L^{\infty}(\Omega)} + d)) d\tilde{t}\right) \\ &\geq \exp\left(\left(\inf_{\varrho \in \mathbb{R}_{+}} \mathcal{Q}(\varrho) - (\|s\|_{L^{\infty}} + d)\right) t\right) =: b(t) > 0 \end{aligned}$$

which is positive definite for every finite time and the bound from below is uniform with respect to δ . Then one remarks that

$$\partial_t \varrho_{\delta} \ge \inf_{x \in \Omega} s(x) \overline{n}_{\delta}(t) \ge \inf_{x \in \Omega} s(x) b(t) > 0$$

which proves that ρ_{δ} is strictly increasing for any fixed time uniformly with respect to δ provided condition (iii)' of Hypotheses 3.2.

– Reaching ρ_0 :

As ρ_{δ} is an increasing function whose initial datum is zero, there are two possibilities

- 1. Either ρ_{δ} never reaches ρ_0 *i.e.* $\rho(t) < \rho_0$ for every non negative $t \in \mathbb{R}$. Then for all times $Q_{\delta}(\rho_{\delta}) = Q(\rho_{\delta})$ and Q is always regular thus uniqueness results from Theorem 1 imply that $n_{\delta} = n$ a.e. in $(0, \infty) \times \Omega$. There is nothing to prove
- 2. Or there exists a time t_0 s.t. $\rho_{\delta}(t_0) = \rho_0$. Again by uniqueness, one has that

$$n_{\delta}(t, x) = n(t, x)$$
 a.e $(t, x) \in (0, t_0) \times \Omega$.

and thus this time t_0 is equal for every δ . We fix a time $t_1 = 2t_0$ and define \underline{c} as

$$\underline{c} := \inf_{t \in [0,t_1]} \overline{n}(0) \exp\left(\left(\inf_{\varrho \in \mathbb{R}_+} \mathcal{Q}(\varrho) - (\|s\|_{L^{\infty}} + d)\right) t\right) > 0$$

We deduce then that in the neighborhood of t_0 one can write

$$\varrho_{\delta}(t_0+\omega)-\varrho_0>\underline{c}\omega, \quad \forall \omega>0.$$

Choosing then $\omega = \delta/\underline{c}$, one writes

$$\begin{split} &\int_{\mathcal{O}_{t_1}} (Q_{\delta}(\varrho_{\delta})n_{\delta} - Q(\varrho)n)\eta dx dt \\ &= \int_{((0,t_0)\cup(t_0+\omega,t_1))\times\Omega} (Q(\varrho_{\delta})n_{\delta} - Q(\varrho)n)\eta dx dt \\ &+ \int_{((t_0,t_0+\omega))\times\Omega} (Q_{\delta}(\varrho_{\delta})n_{\delta} - Q(\varrho)n)\eta dx dt =: R_1 + R_2 \end{split}$$

As previously shown $n_{\delta}(t, x) \equiv n(t, x)$ everywhere on $(0, t_0)$ a.e. $x \in \Omega$, thus R_1 reduces to

$$R_{1} = \int_{((t_{0}+\omega,t_{1}))\times\Omega} (Q(\varrho_{\delta})n_{\delta} - Q(\varrho)n)\eta dxdt$$

$$= \int_{((t_{0}+\omega,t_{1}))\times\Omega} (Q(\varrho_{\delta})n_{\delta} - Q(\varrho)n_{\delta} + Q(\varrho)n_{\delta} - Q(\varrho)n)\eta dxdt$$

$$\leq c_{1} \|\varrho_{\delta} - \varrho\|_{L^{\infty}(0,t_{1})} \|n_{\delta}\|_{L^{\infty}((0,t_{1});L^{2}(\Omega))}$$

$$+ \left| \int_{((t_{0}+\omega,t_{1}))\times\Omega} (n_{\delta} - n)Q(\varrho)\eta dxdt \right|$$

where $c_1 = \|Q\|_{W^{1,\infty}(\mathbb{R}\setminus B(\varrho_0,\delta))}$, the latter term tends to zero when $\delta \to 0$ thanks to weak convergence arguments on n_{δ} . On the other hand,

$$R_2 \le \sqrt{\omega} \|Q\|_{L^{\infty}} \left(\|n\|_{L^{\infty}(0,t_1;L^2(\Omega))} + \|n_{\delta}\|_{L^{\infty}(0,t_1;L^2(\Omega))} \right) \|\eta\|_{H^1(\mathcal{O}_{t_1})}$$

This proves existence of a weak solution $n \in V(\mathcal{O}_T)$ of (1). Uniqueness follows up to the time t_0^- .

We end this part with a time asymptotic result.

Proposition 1 Under hypotheses 3.1 and supposing that

$$d>\lim_{\varrho\to\infty}Q(\varrho),$$

the solution of (1) provides a monotone increasing function $\varrho(t) = \int_0^t \int_{\Omega} s(x)n(\tilde{t}, x) dx d\tilde{t}$ that satisfies

$$\varrho^{\infty} := \lim_{t \to \infty} \varrho(t) < \infty$$

Proof One has

$$\frac{d}{dt}\varrho = \int_{\mathbb{R}} s(x)n(t,x)dx \ge 0,$$

and thus $\rho(t)$ is monotone increasing. By contradiction, assume that $\rho(t) \to \infty$. Setting $\overline{n}(t) := \int_{\Omega} n(t, x) dx$ and testing the weak formulation with 1 gives

$$\left[\overline{n}(t) + \varrho(t)\right]_{t=t_1}^{t=t_2} = \int_{t_1}^{t_2} \left[Q(\varrho(t)) - d \right] \overline{n}(t) dt$$

Since $d > \lim_{t\to\infty} Q(\varrho(t))$, for t_1 large enough, the right hand side becomes negative implying $\overline{n}(t) + \varrho(t)$ shall not increase any more. This contradicts the assumption $\varrho(t) \to \infty$. We conclude that $\varrho^{\infty} < \infty$.

3.2 Spectral analysis

In this part, we investigate the spectral decomposition of the solution, in particular under the assumption of piecewise-linear coefficient Q(.).

3.2.1 Spectral decomposition

As it is usual in the field of parabolic equations (Dautray and Lions 1990), one shall try the variable separation. To this aim we consider the spectral problem in the trait-space Ω .

We denote \mathcal{A} the operator defined by:

$$\mathcal{A}n(x) = -\mu \partial_{xx}^2 n(x) + s(x)n(x).$$

Lemma 1 We suppose that s is a bounded function for almost every $x \in \Omega$. We denote the eigenvalues Λ_k (resp. the eigenvectors V_k) of A the solution of the following equation:

$$\begin{cases} \mathcal{A}V_k = \Lambda_k V_k, & a.e. \ x \in \Omega, \ \forall k \in \mathbb{N} \\ \partial_x V_k(-1) = \partial_x V_k(1) = 0, & x \in \partial \Omega \end{cases}$$
(4)

All eigenvalues are simple. The sequence of eigenvalues $(\Lambda_k)_{k\in\mathbb{N}}$ is monotone increasing and positive. The limit of Λ_k when k goes to infinity is infinite. Moreover the sequence $(V_k)_{k\in\mathbb{N}}$ is an orthonormal basis of $L^2(\Omega)$. One can bound the eigenvalues, setting $m := \inf_{x\in\Omega} \min(\mu, s)$ and $M := \sup_{x\in\Omega} \max(\mu, s)$,

$$m\left(\left(\frac{k\pi}{2}\right)^2+1\right) \leq \Lambda_k \leq M\left(\left(\frac{k\pi}{2}\right)^2+1\right).$$

The eigenvector associated to the smallest eigenvalue is positive. The eigenvector associated to the k^{th} eigenvalue Λ_k has precisely k zeros on Ω .

Proof The proof is standard (Coddington and Levinson 1955; Dautray and Lions 1990) and can be found for instance in Zettl (2005) Theorem 4.6.2 p. 87. □

Then we project the initial data on the eigen-basis. The solution n of (1) shall be expressed as a time dependent superposition of modes:

$$n = \sum_{k \in \mathbb{N}} \alpha_k(t) V_k(x), \quad \forall (t, x) \in \mathbb{R}_+ \times \Omega$$

where the α_k should satisfy

$$\begin{cases} \partial_t \alpha_k + (\Lambda_k + d)\alpha_k = Q\left(\sum_k s_k \int_0^t \alpha_k(\tilde{t})d\tilde{t}\right)\alpha_k, & k \ge 1\\ s_k = \int_{\Omega} V_k(x)s(x)dx = \Lambda_k \int_{\Omega} V_k dx, & (5)\\ \alpha_k(0) = \int_{\Omega} V_k(x)n_I(x)dx \end{cases}$$

We now make the assumption that the non-linearity Q is piecewise-constant:

Assumption 3.1 The function Q(.) is piecewise-constant: $Q(\varrho) = Q_0$ for $\varrho \le \varrho_0$ and $Q(\varrho) = Q_1$ for $\varrho > \varrho_0$.

Under this assumption, and for a time small enough $t < t_0$, *i.e* before $\rho(t)$ reaches ρ_0 , one solves explicitly (5) which gives:

$$\begin{cases} \alpha_k = \alpha_k(0) \exp\left((Q_0 - d - \Lambda_k)t\right), & t \in [0, t_0[\\ \alpha_k(0) = \langle V_k, n_I \rangle \end{cases}$$

where < ., . > denotes the usual scalar product on $L^2(\Omega)$. We denote by $\Lambda_k^0 := b - \Lambda_k$ where $b := Q_0 - d$, so that we can write the solution as:

$$n(t,x) = \sum_{k=0}^{\infty} \langle n_I, V_k \rangle e^{\Lambda_k^0 t} V_k(x), \quad t \le t_{\varrho_0}$$
(6)

Let us define:

$$t_{\rho_0} = \inf\{t \ge 0; \ \rho(t) = \rho_0\}$$
(7)

Then, one estimates the value of t_{ρ_0} , from the above expression, and solve the system for time greater than t_{ρ_0} . To find the value of t_{ρ_0} , one solves the non-linear equation $\rho_0 = \rho(t_{\rho_0})$, more precisely:

$$\varrho_0 = \sum_k \phi_k \left(e^{A_k^0 t_{\varrho_0}} - 1 \right) \tag{8}$$

with

$$\phi_k := \frac{\langle n_I, V_k \rangle \langle V_k, s \rangle}{\Lambda_k^0}.$$

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It appears relatively difficult to solve analytically (8). However, we will study in Sects. 4 and 5 several approaches to derive asymptotic estimates of this key quantity.

It remains to solve the system for time t larger than t_{ϱ_0} . The strategy is exactly the same, except that the new eigenvalues are now $\Lambda_k^1 = (Q_1 - d) - \Lambda_k$, and that the new initial condition is $n(t_{\varrho_0}, x)$, denoted $n_{t_{\varrho_0}}(x)$. Therefore, for $t > t_{\varrho_0}$:

$$n(t, x) = \sum_{k} < n_{t_{\varrho_0}}, V_k > e^{\Lambda_k^1(t - t_{\varrho_0})} V_k(x)$$
(9)

Since the spectrum of A is located on the positive real axis : if $d > Q_1$ then

$$\lim_{t \to \infty} n(t, x) = 0, \quad a.e. \, x \in \Omega$$

3.2.2 Spectral calculus

We now consider the problem of finding explicit expressions for the eigenvalues and eigenvectors associated with problem (4). To this end, we now make some more assumptions on Q, s and μ :

Assumption 3.2 1. The function Q(.) is piecewise-constant: $Q(\varrho) = Q_0$ for $\varrho \le \varrho_0$ and $Q(\varrho) = Q_1$ for $\varrho > \varrho_0$

- 2. The function s(.) is piecewise-constant: $s(x) = s_{\varepsilon}(x) = 1$ if $x \in [-\varepsilon, \varepsilon] \cap \Omega$ and 0 otherwise.
- 3. The function $\mu(.)$ is constant $\mu(x) = \mu > 0$.

Notice that ε is a parameter in (0, 1] that we do not assume, here, to be small. We consider this limit in Sect. 4 instead.

We study the eigen-problem (4) under Assumptions 3.2. This means find the sequence (V_k, Λ_k) solving

$$-\mu V_k'' + s_{\varepsilon} V_k = \Lambda_k V_k$$
, in]0, 1[, $V_k'(0) = V_k'(1) = 0$.

On each part of the domain one has a constant coefficient problem that can be solved. Our goal is to construct by composition the complete eigenproblem in this particular case. Converting the second order problem into a first order system one has:

$$\begin{cases} \partial_x \mathbf{Y}(x, \Lambda) = M(x, \Lambda) \mathbf{Y}(x, \Lambda), & \text{in }]0, 1[\\ Y_2(0, \Lambda) = Y_2(1, \Lambda) = 0, & x \in \{0, 1\} \end{cases}$$
(10)

where $\mathbf{Y} = (V_k, \partial_x V_k)^T$. So the complete problem can be solved by piecewise exponentials:

$$\mathbf{Y}(x,\Lambda) = \boldsymbol{\Phi}(x,\Lambda)\mathbf{Y}(0), \quad \boldsymbol{\Phi}(x,\Lambda) := \begin{cases} e^{M(\Lambda-1)x} & \text{if } x < \varepsilon \\ e^{M(\Lambda)(x-\varepsilon)}e^{M(\Lambda-1)\varepsilon} & \text{if } x > \varepsilon \end{cases},$$

and

$$M(\Lambda) = \begin{pmatrix} 0 & 1 \\ -\omega_0^2(\Lambda) & 0 \end{pmatrix} \mathbb{1}_{(0,\varepsilon)}(x) + \begin{pmatrix} 0 & 1 \\ -\omega_1^2(\Lambda) & 0 \end{pmatrix} \mathbb{1}_{(\varepsilon,1)}(x)$$

where $\omega_0^2 := -\frac{s_0 - \Lambda}{\mu}$ and $\omega_1^2 := \frac{\Lambda}{\mu}$. The exponential matrices are then explicit:

$$e^{M(\Lambda)x} = \begin{pmatrix} \cos(\omega_0 x) & \sin(\omega_0 x)/\omega_0 \\ -\omega_1 \sin(\omega_0 x) & \cos(\omega_0 x) \end{pmatrix} \mathbb{1}_{(0,\varepsilon)}(x) \\ + \begin{pmatrix} \cos(\omega_1 x) & \sin(\omega_1 x)/\omega_1 \\ -\omega_1 \sin(\omega_1 x) & \cos(\omega_1 x) \end{pmatrix} \mathbb{1}_{(\varepsilon,1)}(x).$$

The boundary problem (10) can be rewritten in an algebraic form:

$$T\mathbf{Y}(0) = 0, \quad T := N_1 + N_2 \Phi(1), \quad N_1 := \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix}, \quad N_2 := \begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix},$$

that has a solution when det T = 0 i.e.

$$\left|N_1 + N_2 e^{M(\Lambda)(1-\varepsilon)} e^{M(\Lambda-1)\varepsilon}\right| = 0$$

this leads to a simpler condition

 $\omega_0 \tan_0 + \omega_1 \tan_1 = 0$

where $\tan_j := \tan(\omega_j x_j)$, $j \in \{0, 1\}$, $x_0 = \varepsilon$ and $x_1 = (1 - \varepsilon)$. As a function of Λ the latter equation becomes:

$$\sqrt{\Lambda - 1} \tan\left(\delta\sqrt{\Lambda - 1}\mu\right) + \sqrt{\Lambda} \tan\left((1 - \delta)\sqrt{\Lambda}\right) = 0.$$
(11)

The corresponding eigenvector is

$$V(x) = \mathbb{1}_{x \le \varepsilon} \cos(x\omega_0) + \mathbb{1}_{x \ge \varepsilon} \left\{ \cos(\varepsilon\omega_0) \cos(\omega_1(x-\varepsilon)) - \sin(\varepsilon\omega_0) \sin(\omega_1(x-\varepsilon)) \frac{\omega_0}{\omega_1} \right\},$$
(12)

and one should simply take into account whether ω_0 is a pure imaginary or real number in order to pass from hyperbolic to standard trigonometric functions. So there exists a sequence $(\Lambda_k, V_k)_{k \in \mathbb{N}}$, s.t. for each $k \in \mathbb{N}$, Λ_k solves (11) and V_k writes as in (12). But, for a given set of data s(x) and μ , Λ_k , the solution (11) is not explicit. In what follows we approximate it by a spectral asymptotic expansion.

4 Asymptotic analysis for narrow selection profiles

In this section, we focus our attention on the case of a narrow selection profile, namely considering that the selection function s(x) can be written as $s(x/\epsilon)$ where ϵ is a small parameter. In biological terms, this assumption means that selection is very specific, and that only B-cells with a trait very similar to the target are selected.

We first consider this asymptotic regime from the spectral point of view, and then construct an asymptotic expansion of the solution in ϵ . These two approaches enable an asymptotic estimation of the time to threshold, characterizing the duration of the production process until a sufficient level of selected B-cell is reached.

4.1 Asymptotic spectral analysis

4.1.1 Spectrum of auto-adjoint operators with compact inverse

Because the Neumann problem admits a zero eigenvalue we shift the spectrum by adding the identity because then $-\partial_{x^2}^2 + I$ has an auto-adjoint compact inverse and the spectral theory can be used. Since the shift is artificial and does not change any of the results presented below we omit it—in what follows we return to the original operator. Any time that we mention that the operator has a compact inverse, it is understand in the sense above (i.e., modulo a unit shift).

In a first step we study again problem (4) under assumptions 3.2. Since the operator \mathcal{A} is auto-adjoint and has a compact inverse from $L^2(0, 1)$ into itself $L^2(0, 1)$ it admits (see Theorem 6. p. 38 Dautray and Lions 1990) a discrete spectrum that can be arranged into an increasing sequence of real eigenpairs denoted ($\Lambda_{\varepsilon,k}, V_{\varepsilon,k}$) for $k \in \mathbb{N}$.

4.1.2 An asymptotic Ansatz

We develop $V_{\varepsilon,k}$ using an asymptotic expansion, to this purpose we approach $V_{\varepsilon,k}$ by a series reading

$$\mathcal{V}_{\varepsilon,k}(x, y) := v_{0,k}(x, y) + \varepsilon v_{1,k}(x, y) + \varepsilon^2 \pi_{\varepsilon,2,k}(x, y)$$

where x represents the slow variable $(x \in (0, 1))$ and y the fast variable (typically $y = x/\varepsilon$). Replacing this expansion in the eigen-problem and gathering terms for each power of ε provides equations at each scale. For instance, at the order ε^{-2} one recovers that $v_{0,k}$ does not depend on y, while at order ε^{0} it solves

$$\begin{cases} -\mu v_{0,k}'' = \lambda_{0,k} v_{0,k}, & x \in (0,1), \\ v_{0,k}'(0) = v_{0,k}'(1) = 0, \end{cases}$$
(13)

the solution is explicit : $\lambda_{0,k} := \mu(k\pi)^2$ for $k \in \mathbb{Z}$, whereas normalizing the eigenvectors gives:

$$v_{0,k}(x) = \begin{cases} 1, & \text{if } k = 0\\ \sqrt{2}\cos(k\pi x), & \text{otherwise} \end{cases}$$

Since $v_{0,k}$ does not solve the original equation, we correct it by adding a second order microscopic corrector $\pi_{\varepsilon,2,k}(x) := \pi_2(x/\varepsilon)v_{0,k}(x)$ where $\pi_2(y)$ solves

$$\begin{cases} -\mu \partial_{y^2}^2 \pi_2 + s_0(y) = 0, & x \in (0, \infty), \\ \partial_y \pi_2(0) = 0. \end{cases}$$
(14)

The solution is explicit:

$$\pi_2(y) = \pi_2(0) + \frac{1}{\mu} \int_0^y \int_0^z s(\tilde{z}) d\tilde{z} dz$$

Since all the functions are defined up to a constant we omit them in the rest of the section. If we take the first order derivative of $\pi_{2,k}$ one has:

$$\partial_y \pi_{2,k}(y) = \frac{1}{\mu} \begin{cases} \int_0^y s(z) dz & \text{if } y < 1, \\ \overline{s} := \int_0^1 s(z) dz & \text{otherwise} \end{cases}.$$

In order to reduce the contribution of the growth at infinity of the latter microscopic boundary layer function, we introduce the first order spectral problem: $v_{1,k}$ solves

$$\begin{cases} -\mu v_{1,k}'' = \lambda_{0,k} v_{1,k} + \lambda_{1,k} v_{0,k} + 2\bar{s} v_{0,k}', & x \in (0, 1), \\ v_{1,k}'(0) = 0, & \text{if } x = 0, \\ v_{1,k}'(1) = -\frac{\bar{s}}{\mu} v_{0,k}(1), & \text{if } x = 1. \end{cases}$$

Because the operator $-\partial_{x^2}^2 - \lambda_{0,k}$ admits a kernel of dimension one and is auto-adjoint, the data of this problem should be polar to the kernel of this operator. Therefore, due to this latter condition:

$$\lambda_{1,k} := \frac{\overline{s} \ v_{0,k}^2(0)}{\|v_{0,k}\|_{L^2(0,1)}^2} = \begin{cases} 2\overline{s} & \text{if } k \neq 0, \\ \overline{s} & \text{otherwise.} \end{cases}$$

Interestingly enough this first order eigen-contribution is independent on μ . The solution of the latter problem then reads (modulo a multiple of $v_{0,k}$)

$$v_{1,k} = -\frac{\bar{s}}{\mu} \cdot \begin{cases} x^2/2 & \text{if } k = 0, \\ (1-x)\frac{v'_{0,k}(x)}{(k\pi)^2} + xv_{0,k}(x) & \text{otherwise} \end{cases}$$

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We compute the problem solved by $\mathcal{V}_{\varepsilon,k}$:

$$\begin{split} -\mu \mathcal{V}_{\varepsilon,k}'' + s_{\varepsilon} \mathcal{V}_{\varepsilon,k} &= -\mu v_{0,k}'' + s_{\varepsilon} v_{0,k} - \varepsilon \mu v_{1,k}'' + \varepsilon s_{\varepsilon} v_{1,k} \\ &- \varepsilon^2 \mu \Biggl(\frac{\partial_{y^2}^2 \pi_{2,k}(x/\varepsilon)}{\varepsilon^2} v_{0,k} + 2 \frac{\partial_{y} \pi_{2,k}(x/\varepsilon)}{\varepsilon} v_{0,k}' + \pi_{2,k} v_{0,k}'' \Biggr) \\ &+ \varepsilon^2 s_{\varepsilon} \pi_{\varepsilon,2,k} \\ &= (\lambda_{0,k} + \varepsilon \lambda_{1,k}) (v_{0,k} + \varepsilon v_{1,k} + \varepsilon^2 \pi_{\varepsilon,2,k}) - \varepsilon^2 (\lambda_{1,k} v_{1,k} + \varepsilon \pi_{\varepsilon,2,k}) \\ &+ \varepsilon s_{\varepsilon} (v_{1,k} + \varepsilon \pi_{\varepsilon,2,k}) + 2\varepsilon \left(\int_x^{\varepsilon} s(z/\varepsilon) dz \mathbb{1}_{(0,\varepsilon)}(x) \right) v_{0,k}' \Biggr) \\ &= \lambda_{\varepsilon,1,k} \mathcal{V}_{\varepsilon,k} + \varepsilon \left(s_{\varepsilon} v_{1,k} + 2 \left(\int_x^{\varepsilon} s(z/\varepsilon) dz \mathbb{1}_{(0,\varepsilon)}(x) \right) v_{0,k}' \right) \\ &+ O(\varepsilon^2). \end{split}$$

where we defined $\lambda_{\varepsilon,1,k} := \lambda_{0,k} + \varepsilon \lambda_{1,k}$. We define the bi-continuous bi-linear form associated to the elliptic operator above:

$$a_{\varepsilon}(u,v) := \int_{\Omega} u'(x)v'(x)dx + s_{\varepsilon}(x)u(x)v(x)dx$$

which is coercive on $H^1(\Omega)$ for any positive definite ε .

Theorem 3 Under the hypotheses on s_{ε} , and for ε_0 small enough, the tuple $(\lambda_{\varepsilon,1,k}, \mathcal{V}_{\varepsilon,k})$ is a generalized eigen-pair, i.e., it verifies

$$\left|a_{\varepsilon}(\mathcal{V}_{\varepsilon,k},v)-\lambda_{\varepsilon,1,k}(\mathcal{V}_{\varepsilon,k},v)\right| \leq C_{1}(k+1)\pi\varepsilon^{\frac{3}{2}} \left\|\mathcal{V}_{\varepsilon,k}\right\|_{L^{2}(0,1)} \|v\|_{L^{2}(0,1)}, \quad \forall v \in H^{1}(0,1),$$

which implies that

$$\left|\Lambda_{\varepsilon,k}-\lambda_{\varepsilon,1,k}\right|\leq C_2(k+1)\pi\varepsilon^{\frac{3}{2}}.$$

The constants C_1 and C_2 depend on s, but not on μ .

Proof A simple triangular inequality shows that there exists a constant c_0 independent on ε and on k s.t.

$$0 < c_0 \leq \left\| \mathcal{V}_{\varepsilon,k} \right\|_{L^2(\Omega)}, \quad \forall k \in \mathbb{N}, \ \varepsilon \in (0, \varepsilon_0),$$

for ε_0 small enough. Then $\mathcal{U}_{\varepsilon}$ being regular (typically $\mathcal{U}_{\varepsilon} \in H^2(\Omega)$), the error estimates computed above hold as well in the weak formulation associated and lead to estimate from above:

$$\begin{split} \varepsilon \left(\left(s_{\varepsilon} v_{1,k} + 2 \left(\int_{x}^{\varepsilon} s(z/\varepsilon) dz \mathbb{1}_{(0,\varepsilon)}(x) \right) v_{0,k}' \right), v \right)_{L^{2}(0,1)} \\ &\leq \varepsilon \| s_{\varepsilon} \|_{L^{2}(0,1)} \| v_{1,k} \|_{L^{\infty}(0,1)} \| v \|_{L^{2}(0,1)} \\ &+ 2\varepsilon \| s_{\varepsilon} \|_{L^{\infty}} \| \mathbb{1}_{0,\varepsilon} \|_{L^{2}(0,1)} \| v_{0,k}' \|_{L^{\infty}(0,1)} \| v \|_{L^{2}(0,1)} \\ &\leq c_{1} \varepsilon^{\frac{3}{2}} (1+k) \pi \| \mathcal{V}_{\varepsilon,k} \|_{L^{2}(0,1)} \| v \|_{L^{2}(0,1)} \end{split}$$

while higher order terms are bounded in the $L^{\infty}(0, 1)$ norm with respect to ε and independent on k. Then we apply Theorem 6 p. 38 in Dautray and Lions (1990), whose proof by contradiction, based on the spectral decomposition of a_{ε} , can be found in Bendali et al. (2009) and Bendali et al. (2008).

4.1.3 Numerical comparison

In Fig. 1, we compare the first order approximation of the eigenvalues for the first ten values of k and we display the numerical solution of the characteristic polynomial (11). We choose $\varepsilon = 0.1$.

4.1.4 Application: estimation of the time to threshold

We consider the non-linear equation:



Fig. 1 Numerical comparison between asymptotic and numerical eigen-pairs. On the *left* the eigenvectors for $k \in \{0, 1\}$ and $\varepsilon = 0.1$, on the *right* the eigenvalues for $k \in \{0, ..., 11\}$

with

$$\Phi_{\varepsilon,k} := \frac{\langle n_I, V_{\varepsilon,k} \rangle \langle V_{\varepsilon,k}, s \rangle}{b - \Lambda_{\varepsilon,k}}$$

There exists no analytical solution for t_{ϱ_0} for the above equation. However, it is possible to obtain an asymptotic description of t_{ϱ_0} . In Sect. 3.2, we derived an asymptotic expansion of the eigenvalues $\Lambda_{\varepsilon,k}$ and eigenvectors $V_{\varepsilon,k}$ when $\epsilon \to 0$ in the form:

$$\Lambda_{\varepsilon,k} = \sum_{l \ge 0} \lambda_{l,k} \epsilon^l \quad \text{and} \quad V_{\varepsilon,k} = \sum_{l \ge 0} v_{l,k} \epsilon^l \tag{16}$$

We first show the following:

Lemma 2 Assuming that the following limits exists:

$$\lim_{\varepsilon \to 0} \frac{1}{\varepsilon} \int s(x/\varepsilon) v_{i,k}(x) dx = \underline{v}_{i,k}$$
(17)

and that $\lambda_{0,k} \neq b$ for all $k \geq 0$, the coefficient $\Phi_{\varepsilon,k}$ can be approximated by: $\Phi_{\varepsilon,k} = \varepsilon \phi_{0,k} + O(\varepsilon^2)$, with $\phi_{0,k} = \frac{\langle n_I, v_k^0 \geq v_k^0}{b - \lambda_{0,k}}$.

Proof By definition,

$$\Phi_{\varepsilon,k} := \frac{\langle n_I, v_k^{\epsilon} \rangle \langle s^{\epsilon}, v_k^{\epsilon} \rangle}{b - \Lambda_k^{\epsilon}}$$
(18)

From the asymptotic description of the spectrum, we know that:

$$\Lambda_{\varepsilon,k} = \lambda_{0,k} + \varepsilon \lambda_{1,k} + O(\varepsilon^2) \tag{19}$$

$$V_{\varepsilon,k} = v_{0,k} + \epsilon v_{1,k} + O(\varepsilon^2)$$
⁽²⁰⁾

Therefore, assuming that the following limits exists:

$$\lim_{\epsilon \to 0} \frac{1}{\epsilon} \int s(x/\epsilon) v_{i,k}(x) dx = \underline{v}_{i,k}$$
(21)

we obtain:

$$\langle s_{\epsilon}, V_{\varepsilon,k} \rangle = \varepsilon \underline{v}_{0,k} + \epsilon^2 \underline{v}_{1,k} + O(\epsilon^3)$$
 (22)

Gathering the above estimates, we obtain the desired result.

An asymptotic description of t_{ϱ_0} is given by the following result, which essentially assumes that only the first mode is growing (assumption (iii) below).

Lemma 3 If t_{ϱ_0} is the solution of Eq. (15), with:

(i)
$$\Lambda_{\varepsilon,k} = \lambda_{0,k} + O(\epsilon)$$

(ii) $\Phi_{\epsilon,k} = \epsilon \phi_{0,k} + O(\epsilon^2)$
(iii) $b - \lambda_{0,0} > 0$ and $b - \lambda_{0,k} < 0$ for all $k > 0$

Then t_{ρ_0} diverges to $+\infty$ when $\epsilon \to 0$ as:

$$t_{\varrho_0} = \frac{1}{b - \lambda_{0,0}} \ln\left(\frac{\varrho_0}{\phi_{0,0}} \frac{1}{\epsilon} + o\left(\frac{1}{\epsilon}\right)\right)$$
(23)

Proof We introduce $x = e^{t_{\varrho_0}}$ and look for an expansion of the form:

$$x = \frac{x_0}{\epsilon^{\beta}} + o\left(\frac{1}{\epsilon^{\beta}}\right) \tag{24}$$

First, we write $x^{b-\Lambda_{\varepsilon,k}} = C_{k,\epsilon} \left(\frac{x_0}{\epsilon^{\beta}}\right)^{b-\lambda_{0,k}}$ where $C_{k,\epsilon} \to 1$ when $\epsilon \to 0$. Furthermore, using assumption (iii), one controls the convergence of $x^{b-\Lambda_{\varepsilon,k}}$ to 0 uniformly in *k*. Therefore,

$$\frac{\varrho_0}{\epsilon} = \Phi_{\epsilon,0} C_{0,\epsilon} \left(\frac{x}{\epsilon^{\beta}}\right)^{b-\lambda_{0,0}} - \sum_{k\geq 0} \Phi_{\epsilon,k} + \sum_{k\geq 1} \Phi_{\epsilon,k} x^{b-\Lambda_{\epsilon,k}}$$
(25)

$$=\phi_{0,0}\left(\frac{x_0}{\epsilon^{\beta}}\right)^{b-\lambda_{0,0}} + O(1)$$
(26)

We deduce that

$$\frac{\varrho_0}{\phi_{0,0}x^{b-\lambda_{0,0}}\epsilon} = \frac{1}{\epsilon^{\beta(b-\lambda_{0,0})}}$$
(27)

implying that $\beta = 1/(b - \lambda_{0,0})$ and $x_0 = (\varrho_0/\phi_{0,0})^{1/(b - \lambda_{0,0})}$.

In terms of the original parameters of the model, we conclude that the time to produce an output quantity ρ_0 is asymptotically given by:

$$t_{\varrho_0} = \frac{1}{Q_0 - d} \ln \left(\frac{\varrho_0(Q_0 - d)}{\langle n_I, v_{0,0} \rangle \underline{v}_{0,0}} \frac{1}{\epsilon} + o\left(\frac{1}{\epsilon}\right) \right)$$
(28)

This formula relates in a compact form the birth rate Q_0 , the death rate d, the width of the selection function ϵ , the mutation rate μ and the initial condition n_I to the characteristic time-scale of the B-cell production process. Notice that this formula relies on the assumption that only the first mode grows, meaning that μ must be larger than $(Q_0 - d)/\pi^2$. Numerical comparison between the time t_{Q_0} computed from the numerical solution of the PDE and this formula is displayed in Fig. 2. From the spectral decomposition, we learn that decreasing the mutation rate μ has the effect of recruiting further modes, whereas a large μ implies that the evolution of the solution forgets the other modes which were present in the initial condition.



Fig. 2 Time to threshold t_{ϱ_0} as a function of the threshold parameter ρ_0 in log-scale. The crosses are obtained by a direct numerical simulations of the PDE model with constant birth-rate Q. The line is obtained by evaluating formula (28) at first order, i.e., discarding the terms of $o(1/\epsilon)$. Parameters: b = Q - d = 0.1, $\epsilon = 0.01$

4.2 Asymptotic expansion

4.2.1 Asymptotic expansion of the solution

When Q is piecewise constant, using the notation $b = Q_0 - d$, n(t, x), the solution of (1), solves as well the following linear problem, until $\rho(t)$ reaches ρ_0 :

$$\begin{cases} \partial_t n_{\varepsilon} - \mu \partial_{x^2}^2 n_{\varepsilon} = (b - s_{\varepsilon}) n_{\varepsilon} & (t, x) \in \mathcal{O}_T, \\ \partial_x n_{\varepsilon}(t, x) = 0 & (t, x) \in \Sigma_T, \\ n_{\varepsilon}(0, x) = n_I(x) & (t, x) \in \{0\} \times \Omega \end{cases}$$
(29)

Dropping the balance term between death and birth, we look for an approximation of the solution of:

$$\begin{cases} \partial_t \mathcal{N}_{\varepsilon} - \mu \partial_{x^2}^2 \mathcal{N}_{\varepsilon} + s_{\varepsilon} \mathcal{N}_{\varepsilon} = 0, & (t, x) \in \mathcal{O}_T, \\ \partial_x \mathcal{N}_{\varepsilon} = 0, & (t, x) \in \Sigma_T, \\ \mathcal{N}_{\varepsilon}(0, x) = n_I(x), & (t, x) \in \{0\} \times \Omega. \end{cases}$$
(30)

We make the Ansatz:

$$\mathcal{N}_{\varepsilon} = \mathcal{N}_0\left(t, x, \frac{x}{\varepsilon}\right) + \varepsilon \mathcal{N}_1\left(t, x, \frac{x}{\varepsilon}\right) + \cdots$$

When we plug it in (30), it gives after separating orders of ε that $\mathcal{N}_0(t, x, y)$ is independent of the fast variable y. Moreover, it solves the homogeneous equations:

$$\begin{cases} \partial_t \mathcal{N}_0 - \mu \partial_{x^2}^2 \mathcal{N}_0 = 0, & (t, x) \in \mathcal{O}_T \\ \partial_x \mathcal{N}_0 = 0, & (t, x) \in \Sigma_T \\ \mathcal{N}_0(0, x) = n_I(x)(t, x) \in \{0\} \times \Omega, \end{cases}$$
(31)

and thus $\mathcal{N}_0(t, x) = \sum_{k \in \mathbb{N}} \overline{n}_I^k \exp(-\lambda_{0,k} t) v_k(x).$

Proposition 2 If $n_1 \in L^2(\Omega) \cap L^{\infty}(\Omega)$ then at zeroth order, one can approach $n_{\varepsilon}(t, x)$ solving (29) by $\exp(bt)\mathcal{N}_0(t, x)$ and the error in the $L^{\infty}((0, T); L^2(\Omega))$ norm is estimated as:

$$\sup_{t\in(0,T)} \|n_{\varepsilon}(t,\cdot) - \exp(bt)\mathcal{N}_0(t,\cdot)\|_{L^2(\Omega)} \le \varepsilon \exp(bT)T\|n_I\|_{L^{\infty}(\Omega)}.$$

Proof One defines the zero order error:

$$\mathcal{E}_0(t,x) := \mathcal{N}_\varepsilon(t,x) - \mathcal{N}_0(t,x)$$

which solves in the strong sense

$$\begin{cases} \partial_t \mathcal{E}_0 - \mu \partial_{x^2}^2 \mathcal{E}_0 + s \mathcal{E}_0 = -s \mathcal{N}_0, & (t, x) \in \mathcal{O}_T \\ \partial_x \mathcal{E}_0 = 0, & (t, x) \in \Sigma_T \\ \mathcal{E}_0(0, x) = 0, & (t, x) \in \{0\} \times \Omega, \end{cases}$$

using standard a priori estimates, one obtains that

$$\frac{1}{2}\partial_t \|\mathcal{E}_0(t,\cdot)\|_{L^2(\Omega)}^2 \le \|s\mathcal{N}_0\|_{L^2(\Omega)} \|\mathcal{E}_0\|_{L^2(\Omega)}$$

which then by dividing both sides by $\sqrt{\|\mathcal{E}_0\|_{L^2(\Omega)}^2 + \delta}$ one gets:

$$\frac{1}{2}\partial_t \sqrt{\left\|\mathcal{E}_0(t,\cdot)\right\|_{L^2(\Omega)}^2 + \delta} \le \left\|s\mathcal{N}_0\right\|_{L^2(\Omega)}.$$

This, integrated in time, provides after passing to the limit $\delta \rightarrow 0$,

$$\|\mathcal{E}_{0}(t,\cdot)\|_{L^{2}(\Omega)} \leq \int_{0}^{t} \|s(\cdot)\mathcal{N}_{0}(\tilde{t},\cdot)\|_{L^{2}(\Omega)} d\tilde{t} \leq \varepsilon t \|\mathcal{N}_{0}\|_{L^{\infty}((0,t)\times\Omega)} \leq \|n_{I}\|_{L^{\infty}(\Omega)} \varepsilon t$$

where we used the maximum principle in order to provide the L^{∞} bound on \mathcal{N}_0 and the result follows.

4.2.2 Estimation of the time to threshold

Using asymptotic expansions above, we estimate the time t_{ρ_0} :

Theorem 4 If s satisfies hypotheses 3.2, and $\mu > b/\pi^2$, then for every given ϱ_0 , there exists a time

$$t_{\varrho_0} := \frac{1}{b} \ln \left(1 + \frac{\varrho_0 b}{\varepsilon \overline{n}_0} \right), \tag{32}$$

and a constant C independent on t_{ρ_0} s.t.

$$|\varrho_0 - \varrho_{\text{out}}(t_{\varrho_0})| \le |\ln(\varepsilon)|\sqrt{\varepsilon}C(\mu, n_I)$$

Proof The previous proposition allows to compute ρ_{app} , an approximation of ρ_{out} which reads:

$$\begin{split} \varrho_{app}(t) &:= \int_0^t \int_{\Omega} s(x) \mathcal{N}_0(t, x) \exp(b\tilde{t}) dx d\tilde{t} \\ &= \frac{(\exp(bt) - 1)}{b} \int_{\Omega} n_I(x) dx \int_{\Omega} s(\tilde{x}) d\tilde{x} + \sum_{k \in \mathbb{N}^*} \frac{(\exp((b - \lambda_k)t) - 1)}{(b - \lambda_k)} \\ &\times < n_I, v_k > < s, v_k > \end{split}$$

where $(\lambda_k)_{k \in \mathbb{N}}$ denote eigenvalues associated to the homogeneous heat equation (31). Under hypotheses 3.2, the previous computation gives

$$\varrho_{\text{app}}(t) := \varepsilon \frac{(\exp(bt) - 1)}{b} \overline{n}_I^0 + \sum_{k \in \mathbb{N}^*} \frac{(\exp((b - \lambda_k)t) - 1)}{(b - \lambda_k)} < n_I, v_k > < s, v_k > \varepsilon$$

Using Cauchy-Schwartz, one has that

$$\begin{aligned} \left| \varrho_{\text{out}}(t) - \varrho_{\text{app}}(t) \right| &\leq \int_{0}^{t} \int_{\Omega} s |n(t, x) - \exp(b t) \mathcal{N}_{0}(\tilde{t}, x)| dx d\tilde{t} \\ &\leq \int_{0}^{t} \|s\|_{L^{2}(\omega)} \|n(t, \cdot) - \exp(b\tilde{t}) \mathcal{N}_{0}(\tilde{t}, \cdot)\|_{L^{2}(\Omega)} d\tilde{t} \leq \varepsilon^{\frac{3}{2}} \|n_{I}\|_{L^{\infty}(\Omega)} \\ &\times \int_{0}^{t} \exp(b\tilde{t}) \tilde{t} d\tilde{t} \leq \varepsilon^{\frac{3}{2}} \|n_{I}\|_{L^{\infty}(\Omega)} \frac{t \exp(bt)}{b} \end{aligned}$$

On the other hand using the explicit expression of ρ_{app} one writes:

$$\left| \varrho_{app} - \left(\frac{\exp(bt) - 1}{b} \right) \overline{n}_{I}^{0} \overline{s} \right| \leq S(t) \sum_{k \in \mathbb{N}^{*}} |\langle n, v_{k} \rangle \langle s, v_{k} \rangle|$$

$$\leq S(t) \|s\|_{L^{2}(\Omega)} \|n_{I}\|_{L^{2}(\Omega)} \leq S(t) \sqrt{\varepsilon} \|n_{I}\|_{L^{2}(\Omega)}$$
(33)

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where we denoted $S(t) := \sup_{k \in \mathbb{N}^*} \left| \frac{\exp((b-\lambda_k)t) - 1}{(b-\lambda_k)} \right|$. Then two cases occur:

(i) either there exists k_0 s.t. $\lambda_{k_0} \leq b \leq \lambda_{k_0+1}$ and then

$$S(t) \leq \frac{\exp((b-\mu\pi^2)t)-1}{b-\lambda_{k_0}},$$

(ii) or $b < \lambda_1 = \mu \pi^2$ and thus

$$S(t) \le \frac{2}{(\lambda_1 - b)} = \frac{2}{(\mu \pi^2 - b)}$$

which justifies our hypothesis between b and μ .

Using a triangular inequality gives:

$$\left|\varrho_{\text{out}}(t_{\varrho_0}) - \varrho_0\right| \le \left|\varrho_{\text{out}}(t_{\varrho_0}) - \varrho_{\text{app}}(t_{\varrho_0})\right| + \left|\varrho_{\text{app}}(t_{\varrho_0}) - \varrho_0\right|$$

which, because of the estimates ii) above, gives

$$|\varrho_{\text{out}}(t_{\varrho_0}) - \varrho_0| \le C \left(\varepsilon^{\frac{3}{2}} t_{\varrho_0} \exp(bt_{\varrho_0}) + \sqrt{\varepsilon} \right) \le C \sqrt{\varepsilon} |\ln(\varepsilon)|.$$

Remark 1 In fact, if *b* is large enough [case (i) above], the previous estimate does not hold since in this case $S(t_{\varrho_0})$ blows as $1/\varepsilon$.

Remark 2 Due to (33), one can not, to our knowledge, improve the accuracy when computing t_{ϱ_0} by increasing the order of the asymptotic expansion since the major source of error comes from this step, when integrating in the selection window the zero order term.

4.2.3 Numerical simulation

The parameters we choose are $\rho_0 = 100$, b = 2 together with a random initial condition n_I . We compute the error between $t_{\rho_0}^{h,k}$, the time to reach ρ_0 and the theoretical formula (32). The direct simulation is made using a P1-Finite Element Method with first order implicit Euler scheme for the time discretization. The same numerical tools are used for the rest of the paper. Results are displayed for various (small) values of ε . We plot, in Fig. 3, the error estimates for two values of μ corresponding to $\mu = b/((k + 1/2)^2 \pi^2$ for $k \in \{0, 1\}$. When k = 0 we are in the hypotheses of the latter theorem, whereas for k = 1 the theoretical error is not uniform with respect to ε . For this specific test-case, the error is comparable in both cases, although greater as μ becomes smaller as expected. The numerical order of convergence is greater than what is predicted theoretically.

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Fig. 3 Starting from a random initial condition, the difference between the numerical $t_{\varrho_0}^{h,k}$ and the limit time to reach t_{ϱ_0} for various values of ε

5 Asymptotic analysis for small and large mutation rates

In this section, we focus on the regimes of small and large mutation rates. We first consider the case of initial conditions restricted to a single trait (Dirac initial data), deriving explicit solutions in this case, and providing asymptotics for the time to threshold. Then, we establish asymptotic expansions of the solution in the regimes $\mu \ll 1$ and $\mu \gg 1$.

5.1 The initial condition is a Dirac mass

5.1.1 The case of the whole space

When the size of the domain goes to infinity we face the problem

$$\begin{cases} \partial_t \mathcal{G}_{\varepsilon} - \mu \partial_{x^2}^2 \mathcal{G}_{\varepsilon} + s_{\varepsilon} \mathcal{G}_{\varepsilon} = b \mathcal{G}_{\varepsilon}, & (t, x) \in \mathbb{R}_+ \times \mathbb{R}, \\ \mathcal{G}_{\varepsilon}(0, x) = \delta_z(x), & x \in \mathbb{R}, \end{cases}$$
(34)

Definition 5.1 We define a very weak solution of (34) the function $\mathcal{G}_{\varepsilon} \in L^2(0, T \times \mathbb{R})$ that solves:

$$\int_{(0,T)\times\mathbb{R}} \mathcal{G}_{\varepsilon} \left(-\partial_t - \partial_{x^2}^2 + (s_{\varepsilon}(x) - b) \right) \varphi dx dt - \varphi(0, z) = 0$$

for every $\varphi \in C(0, T; H^1(\mathbb{R})) \cap L^2(0, T; H^2(\mathbb{R}))$ s.t. $\varphi(T, x) = 0, \forall x \in (0, 1).$

Theorem 5 Providing $s_{\varepsilon} \in L^{\infty}(\mathbb{R})$ there exists a unique very weak function $\mathcal{G}_{\varepsilon} \in L^{2}((0, T) \times \mathbb{R})$ solving (34). Moreover one has the comparison principle:

$$\tilde{\mathcal{G}}_{b-s_{\infty}}(t,x,z) \leq \mathcal{G}_{\varepsilon}(t,x,z) \leq \tilde{\mathcal{G}}_{b}(t,x,z), \quad a.e(t,x) \in (0,T) \times \mathbb{R}$$

where $\tilde{\mathcal{G}}_b$ is the fundamental solution of the heat equation and reads :

$$\tilde{\mathcal{G}}_b(t, x, z) = \frac{1}{\sqrt{4\mu\pi t}} \exp\left(bt - \frac{(x-z)^2}{4\mu t}\right),$$

and $s_{\infty} := \|s\|_{L^{\infty}(\mathbb{R})}$.

Proof The proof of existence and uniqueness relies on duality arguments and the Riesz Theorem and is left to the reader for sake of conciseness. We detail the comparison principle. The maximum principle shows that if we define the backward equation :

$$\begin{cases} -\partial_t \varphi - \partial_{x^2}^2 \varphi + \mathfrak{b}\varphi = \Psi, & (t, x) \in \mathcal{O}_T \\ \partial_x \varphi = 0, & (t, x) \in \Sigma_T \\ \varphi(T, 0) = 0, & (t, x) \in \{T\} \times \Omega \end{cases}$$
(35)

where b is a bounded function in $\mathcal{O}_T, \Psi \in \mathcal{D}(\mathcal{O}_T)$ and $\Psi(t, x) \ge 0$ for all $(t, x) \in \mathcal{O}_T$, then $\varphi(t, x) \ge 0$. Indeed, testing the forward expression of the previous equation in the weak form with $\varphi_h^- := \min(\varphi_h, 0)$ where φ_h is a Steklov approximation of φ i.e.,

$$\varphi_h(t,x) = \frac{1}{h} \int_t^{t+h} \varphi(\tau,x) d\tau$$

and passing then to the limit with respect to h, one has

$$\frac{1}{2} \left[\left\| \varphi^{-}(\tau, \cdot) \right\|_{L^{2}(\Omega)}^{2} \right]_{\tau=0}^{\tau=T} + \int_{\mathcal{O}_{T}} \mathfrak{b}^{+} \varphi \varphi^{-} dx dt \leq - \int_{\mathcal{O}_{T}} \mathfrak{b}^{-} \varphi \varphi^{-} dx dt = - \int_{\mathcal{O}_{T}} \mathfrak{b}^{-} \varphi \varphi^{-} \varphi \varphi^{-} dx dt = - \int_{\mathcal{O}_{T}} \mathfrak{$$

by Gronwall, the result comes easily.

The difference between $\mathcal{G}_{\varepsilon}$ and $\tilde{\mathcal{G}}_{b-s_{\infty}}$ that we denote $\hat{\mathcal{G}} := \mathcal{G}_{\varepsilon} - \tilde{\mathcal{G}}_{b-s_{\infty}}$ solves

$$\int_{\mathcal{O}_T} \hat{\mathcal{G}} \left\{ -\partial_t \varphi - \partial_{x^2}^2 \varphi + (s_\varepsilon - b)\varphi \right\} dx dt = \int_{\mathcal{O}_T} \left\{ s_\infty - s_\varepsilon \right\} \tilde{\mathcal{G}}_{b-s_\infty}(t, x)\varphi(t, x) dx dt$$

As $\tilde{\mathcal{G}}_{b-s_{\infty}}(t, x)$ is non-negative for a.e. $(t, x) \in (0, T) \times \Omega$, and choosing $\mathfrak{b} := (s_{\varepsilon} - b)$ in the dual problem (35), with $\Psi \in \mathcal{D}(\mathcal{O}_T)$ and $\Psi \ge 0$, proves that $\hat{\mathcal{G}}$ is positive almost everywhere in \mathcal{O}_T .

Lemma 4 Suppose that $a \in \mathbb{R}_+$ and set $J_a(t) := \int_0^t \exp\left(at - \frac{1}{t}\right) \frac{d\tilde{t}}{\sqrt{\tilde{t}}}$ then

$$J_a(t) \ge \frac{1}{2e} \exp\left(-\frac{2}{t}\right), \quad \forall t \ge 0.$$

Proof Using Jensen's inequality one writes

$$J_a(t) \ge \int_0^t \exp\left(-\frac{1}{\tilde{t}}\right) \frac{d\tilde{t}}{\sqrt{\tilde{t}}} = \int_{1/t}^\infty \exp(-z) \frac{dz}{z^{\frac{3}{2}}} \ge \left(\int_{1/t}^\infty \exp(-z) dz\right)$$
$$\times \left(\frac{\int_{1/t}^\infty \exp(-z) z dz}{\int_{1/t}^\infty \exp(-z) dz}\right)^{-\frac{3}{2}}$$
$$= \exp\left(-\frac{1}{t}\right) \left(\frac{t}{t+1}\right)^{\frac{3}{2}} \ge \exp\left(-\frac{1}{t}\right) \frac{1}{2} \exp\left(-\frac{t+1}{t}\right) = \frac{1}{2e} \exp\left(-\frac{2}{t}\right)$$

The last inequality comes when estimating $t^{\frac{3}{2}}$ by an exponential, *i.e.* $\forall t \geq 0$, $\exp\left(-\frac{1}{t}\right) < 2t^{\frac{3}{2}}$.

Theorem 6 We suppose that s_{ε} is the characteristic function of the set $(-\varepsilon, \varepsilon)$ and that $z > \varepsilon$. If we denote the time t_{ϱ_0} s.t. $\rho_{\varepsilon}(t) := \int_0^t \int_{\mathbb{R}} s(x) \mathcal{G}_{\varepsilon}(t, x, z) dx dt$ reaches ϱ_0 , one has that $t_{\varrho_0} \to \infty$ when either $\mu \to 0$ or $\mu \to \infty$. Moreover t_{ϱ_0} does not grow faster than any polynomial with respect to μ .

Proof We shall provide a lower bound for $I_{\ell}(t) := 2\varepsilon \int_0^t \tilde{\mathcal{G}}_{b-s_{\infty}}(t, -\varepsilon, z)dt$ and un upper bound for $I_u(t) := 2\varepsilon \int_0^t \tilde{\mathcal{G}}_b(t, \varepsilon, z)dt$. If there exists a time t_u (resp. t_{ℓ}) s.t. $I_{\ell}(t_u) \ge \varrho_0$ (resp. $I_u(t_{\ell}) \le \varrho_0$) then $t_{\varrho_0} \le t_u$ (resp. $t_{\varrho_0} \ge t_{\ell}$).

Because the lower bound from Lemma 4 is not accurate for large values of t, we distinguish two cases:

- Either $t < t_0 := (z + \varepsilon)/(2\sqrt{(b + s_{\infty})\mu})$ and one uses Lemma 4:

$$I_{\ell}(t) \ge \frac{\varepsilon(z+\varepsilon)}{4\mu e\sqrt{\pi}} \exp\left(-\frac{(z+\varepsilon)^2}{2\mu t}\right) =: j_1(t).$$

- Or $t > t_0$ and we split the integral I_ℓ in two parts : integrating first on $(0, t_0)$ and then on (t_0, t) . The first part is estimated thanks to Lemma 4 again, whereas for the latter integral one writes:

$$\begin{split} I_{\ell}(t) &\geq j_{1}(t_{0}) + \frac{\exp(-(b-s_{\infty})t_{0})}{\sqrt{\mu\pi}} \int_{t_{0}}^{t} \exp((b-s_{\infty})\tilde{t}) \frac{d\tilde{t}}{\sqrt{\tilde{t}}} \\ &\geq j_{1}(t_{0}) + \frac{\exp(-(b-s_{\infty})t_{0})}{\sqrt{\mu\pi}k!} \int_{t_{0}}^{t} (b-s_{\infty})^{k} \tilde{t}^{k-\frac{1}{2}} d\tilde{t} \\ &= j_{1}(t_{0}) + \frac{\exp(-(b-s_{\infty})t_{0})}{\sqrt{\mu\pi}k!(k+\frac{1}{2})} (b-s_{\infty})^{k} \left[\tilde{t}^{k+\frac{1}{2}}\right]_{\tilde{t}=t_{0}}^{\tilde{t}=t} =: j_{2}(t). \end{split}$$

since for larger times the linear term $(b - s_{\infty}) t$ dominates $(z + \varepsilon)^2/(4\mu t)$ inside the exponential.

Inverting the latter piecewise formula, one obtains

$$t_{u} := \begin{cases} j_{1}^{-1}(\varrho_{0}) & \text{if } \varrho_{0} < j(t_{0}), \\ j_{2}^{-1}(\varrho_{0}) & \text{otherwise }, \end{cases}$$

where

$$j_1^{-1}(\varrho_0) := \frac{(z+\varepsilon)^2}{-2\ln\left(\frac{\varrho_0 4e\mu}{\varepsilon(z+\varepsilon)}\right)}, \text{ and } j_2^{-1}(\varrho_0)$$
$$:= \left(\frac{(k+\frac{1}{2})k!\sqrt{\pi\mu}}{(b-s_\infty)^k}\exp((b-s_\infty)t_0)(\varrho_0-j_1(t_0)) + t_0^{k+\frac{1}{2}}\right)^{\frac{1}{k+\frac{1}{2}}}$$

This proves the second part of the claim.

In order to provide a lower bound t_{ℓ} , we estimate ρ_{out} by I_u . Indeed we write:

$$I_u(t) \le \frac{2\varepsilon}{\sqrt{\mu\pi}} \exp\left(bt - \frac{(z-\varepsilon)^2}{4\mu t}\right) \sqrt{t} \le \frac{2\varepsilon}{\sqrt{\mu\pi}} \exp\left((b+1)t - \frac{(z-\varepsilon)^2}{4\mu t}\right)$$

this gives that

$$t_{\ell} := \frac{\ln \omega + \sqrt{\ln^2 \omega + (z - \varepsilon)^2 (b + 1)/\mu}}{2(b + 1)}, \quad \omega := \left(\varrho_0 / 2\varepsilon \sqrt{\pi}\right) \sqrt{\mu}$$

which proves the first part of the claim.

In the previous proof the results remain identical if $z < -\varepsilon$ by symmetry. We plot in Fig. 4, left, the comparison between our bounds and numerical computations of the integrals I_l and I_u for a given set of data $(b, \varrho_0, \varepsilon, z)$ when μ varies : t_{ϱ_0} is always between the t_{I_u} and t_{I_l} curves.

5.1.2 The case of a bounded domain

We consider the problem: find $\mathcal{G}_{\varepsilon}$ solution of

$$\begin{cases} \partial_t \mathcal{G}_{\varepsilon} - \mu \partial_{x^2}^2 \mathcal{G}_{\varepsilon} + s_{\varepsilon} \mathcal{G}_{\varepsilon} = b \mathcal{G}_{\varepsilon}, & (t, x) \in \mathbb{R}_+ \times (0, 1), \\ \mathcal{G}_{\varepsilon}'(t, 0) = \mathcal{G}_{\varepsilon}'(t, 1) = 0, & t \in \mathbb{R}_+, \\ \mathcal{G}_{\varepsilon}(0, x) = \delta_z(x), & x \in (0, 1), \end{cases}$$
(36)

where the support of the dirac mass is located in $z \in (0, 1)$.

Theorem 7 If $s \in L^{\infty}(0, 1)$ and $b \in \mathbb{R}$, there exists a unique very weak solution $\mathcal{G}_{\varepsilon} \in L^{2}(\mathcal{O}_{T})$ for every given positive T i.e.,

$$\int_{(0,T)\times\mathbb{R}} \mathcal{G}_{\varepsilon} \left(-\partial_t - \partial_{x^2}^2 + (s(x) - b) \right) \varphi dx dt - \varphi(0, z) = 0$$

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for every test function $\varphi \in C([0, T]; H^1((0, 1))) \cap L^2((0, T); H^2((0, 1)))$. Moreover, one has for a.e. $(t, x) \in \mathcal{O}_T$:

$$\tilde{\mathcal{G}}_{b-s_{\infty}}(t,x,z) \leq \mathcal{G}_{\varepsilon}(t,x,z) \leq \tilde{\mathcal{G}}_{b}(t,x,z),$$

the bounds being fundamental solutions of the heat equation with Neumann boundary conditions:

$$\begin{aligned} \partial_t \tilde{\mathcal{G}}_b - \mu \partial_{x^2}^2 \tilde{\mathcal{G}}_b &= b \tilde{\mathcal{G}}_b, \quad (t, x) \in \mathbb{R}_+ \times (0, 1), \\ \tilde{\mathcal{G}}'_b(t, 0) &= \tilde{\mathcal{G}}'_b(t, 1) = 0, \quad t \in \mathbb{R}_+, \\ \tilde{\mathcal{G}}_b(0, x) &= \delta_z(x), \quad x \in (0, 1), \end{aligned} \tag{37}$$

given explicitly as series by the formula:

$$\tilde{\mathcal{G}}_b(t,x,z) := \frac{\exp(bt)}{\sqrt{4\pi\,\mu t}} \sum_{n\in\mathbb{Z}} \left\{ \exp\left(-\frac{(x-2n+z)^2}{4\mu t}\right) + \exp\left(-\frac{(x-2n-z)^2}{4\mu t}\right) \right\}.$$

We need here two technical lemmas:

Lemma 5 If $x \in (0, \varepsilon)$ and $z \in (\varepsilon, 1)$ and $t \ge 0$ then

$$\tilde{\mathcal{G}}_b(t, x, z) \le \exp\left(bt - \frac{(z-\varepsilon)^2}{4\mu t}\right) \left(\frac{4}{\sqrt{4\pi\mu t}} + \frac{2}{\sqrt{2(1-\varepsilon)}}\right).$$

The proof uses an upper bound of the higher order terms $(|n| \ge 1)$ in the expansion above, and is skipped for sake of conciseness.

Lemma 6 Setting

$$\tilde{\rho}_b(t,z) := \int_0^t \int_0^\varepsilon \tilde{\mathcal{G}}_b(\tilde{t},x,z) dx d\tilde{t},$$

we fix $\varepsilon > 0$, $z \in (\varepsilon, 1)$,

$$\tilde{\rho}_b(t,z) \le \varepsilon \frac{\exp\left((b+1)t - \frac{(z-\varepsilon)^2}{4\mu t}\right)}{\sqrt{4\mu\pi}} \left(1 + \sqrt{\frac{2\mu}{(1-\varepsilon)}}\right),$$

the time that $\tilde{\rho}_b$ reaches ϱ_0 is then greater than

$$t_0 := \frac{\ln \omega + \sqrt{\ln^2 \omega + \frac{(z-\varepsilon)^2}{\mu}}}{2(b+1)}, \text{ and } \omega := \frac{2\varrho_0 \sqrt{(1-\varepsilon)\mu\pi}}{\varepsilon(\sqrt{1-\varepsilon} + \sqrt{2\mu})}.$$

If μ is small then

$$t_0 \sim \frac{|z-\varepsilon|}{2(b+1)} \frac{1}{\sqrt{\mu}}.$$

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When $\mu \to \infty$,

$$t_0 \to \frac{1}{(b+1)} \ln\left(\frac{\varrho_0 \sqrt{2\pi(1-\varepsilon)}}{\varepsilon}\right).$$

Proof Thanks to the previous results there exists a t_1 s.t. $\tilde{\rho}_b(t_1) > \varrho_0$. This implies that t_b , the time such that $\tilde{\rho}_b(t) > \varrho_0$ for all $t > t_b$ is greater than t_1 . As in turn $t_0 \ge t_b$ this ends the proof.

Lemma 7 If $(x, z) \in (0, 1)^2$, the fundamental solution can be estimated from below *as:*

$$\tilde{\mathcal{G}}_b(t,x,z) \ge \exp\left(bt - \frac{(x+z)^2}{4\mu t}\right) \left\{ \frac{1}{\sqrt{4\pi\mu t}} + \frac{1}{4\sqrt{3}} \operatorname{erfc}\left(\sqrt{\frac{3}{\mu t}}\right) \right\}.$$

Proof We write simply that:

$$\sum_{|n|\geq 1} \exp\left(bt - \frac{(x+z-2n)^2}{4\mu t}\right) \geq \exp\left(bt - \frac{(x+z)^2}{4\mu t}\right)$$
$$\times \sum_{|n|\geq 1} \exp\left(-\frac{(x+z)|n|}{\mu t} - \frac{n^2}{\mu t}\right)$$
$$\geq 2\exp\left(bt - \frac{(x+z)^2}{4\mu t}\right) \sum_{n\geq 1} \exp\left(-\frac{3n^2}{\mu t}\right) \geq 2\exp\left(bt - \frac{(x+z)^2}{4\mu t}\right)$$
$$\times \int_{1}^{\infty} \exp\left(-\frac{3s^2}{\mu t}\right) ds$$

which ends the proof.

Again we can estimate the time t_{ϱ_0} for which the selected population ϱ_{out} reaches the threshold value ϱ_0 .

Theorem 8 We suppose that s_{ε} is the characteristic function of the set $(0, \varepsilon)$ and that $z > \varepsilon$, $b \in \mathbb{R}_+$ s.t. $b > s_{\infty} := \|s\|_{L^{\infty}(0,1)}$ and $z > \varepsilon$, then if we denote t_{ϱ_0} the time s.t. ϱ_{out} reaches ϱ_0 , for large (resp. small) values of μ one has:

$$\lim_{\mu \to 0} t_{\varrho_0} = +\infty, \quad \left(\text{resp.} \lim_{\mu \to \infty} t_{\varrho_0} \in (\underline{t}, \overline{t}) \right),$$

where the interval $(\underline{t}, \overline{t})$ depends only on the data set $(\varrho_0, \varepsilon, z, b, s_\infty)$, and $0 < \underline{t} < \overline{t} < \infty$.

Proof When $x < \varepsilon < z$, the heat kernel $\tilde{\mathcal{G}}_b$ is monotone with respect to x. Firstly, we compute the lower bound t_ℓ which is provided estimating $I_u(t) := \varepsilon \int_0^t \tilde{\mathcal{G}}_b(t, \varepsilon, z) dt$

from above. By Lemma 5, one has that

$$I_{u}(t) \leq \varepsilon \frac{\exp\left((b+1)t - \frac{(z-\varepsilon)^{2}}{4\mu t}\right)}{\sqrt{4\mu\pi}} \left(1 + \sqrt{\frac{2\mu}{(1-\varepsilon)}}\right)$$

and so the time that I_u reaches ρ_0 is then greater than

$$t_{\ell} := \frac{\ln \omega + \sqrt{\ln^2 \omega + \frac{(z-\varepsilon)^2}{\mu}}}{2(b+1)}, \text{ and } \omega := \frac{2\varrho_0 \sqrt{(1-\varepsilon)\mu\pi}}{\varepsilon(\sqrt{1-\varepsilon} + \sqrt{2\mu})}$$

If μ is small then

$$t_\ell \sim \frac{|z-\varepsilon|}{2(b+1)} \frac{1}{\sqrt{\mu}}.$$

Whereas if $\mu \to \infty$,

$$t_{\ell} \to \frac{1}{(b+1)} \ln\left(\frac{\varrho_0 \sqrt{2\pi(1-\varepsilon)}}{\varepsilon}\right) =: \underline{t}$$

For what concerns t_u , one has, thanks to Lemma 7,

$$\begin{split} \tilde{\mathcal{G}}_{b-s_{\infty}}(t,x,z) &\geq \tilde{\mathcal{G}}_{b-s_{\infty}}(t,0,z) \geq \exp\left((b-s_{\infty})t - \frac{z^2}{4\mu t}\right) \\ &\left\{\frac{1}{\sqrt{4\pi\,\mu t}} + \frac{1}{4\sqrt{3}}\mathrm{erfc}\left(\sqrt{\frac{3}{\mu t}}\right)\right\}, \end{split}$$

which allows to write:

$$I_{\ell}(t) := \varepsilon \int_0^t \tilde{\mathcal{G}}_{b-s_{\infty}}(t,0,z) dt \ge I_1(t) + I_2(t),$$

where

$$I_1(t) := \frac{\varepsilon}{2\sqrt{\pi\mu}} \int_0^{\tilde{t}} \exp\left((b - s_\infty)\tilde{t} - \frac{z^2}{4\mu\tilde{t}}\right) \frac{d\tilde{t}}{\sqrt{\tilde{t}}},$$

$$I_2(t) := \frac{\varepsilon}{4\sqrt{3}} \int_0^t \exp\left((b - s_\infty)\tilde{t} - \frac{z^2}{4\mu\tilde{t}}\right) \operatorname{erfc}\left(\sqrt{\frac{3}{\mu\tilde{t}}}\right) d\tilde{t},$$

and both functions are increasing with respect to t. This allows us to estimate I_{ℓ} as

$$I_{\ell}(t) \ge I_{1}(t)\mathbb{1}_{(0,t_{0})}(t) + \{I_{1}(t_{0}) + I_{2}(t) - I_{2}(t_{0})\}\mathbb{1}_{(t_{0},\infty)}(t)$$

where, as in the unbounded case, t_0 represents the time at which the coefficient of the exponential function changes sign *i.e.* $t_0 := z/(2\sqrt{\mu(b-s_\infty)})$. As before, when $t < t_0$, the estimate follows the same from Lemma 4. Thus we re-define $j_1(t) := \varepsilon z/(4\mu e) \exp(-z^2/(2\mu t))$. Instead when $t > t_0$, one writes:

$$I_{\ell}(t) \ge j_1(t_0) + \frac{\varepsilon}{4\sqrt{3}} \int_{t_0}^t \exp((b - s_{\infty})(t - t_0)) \operatorname{erfc}\left(\sqrt{\frac{3}{\mu \tilde{t}}}\right) d\tilde{t},$$

but because $\operatorname{erfc}(\sqrt{3/(\mu t)})$ is a monotone increasing function with respect to *t*, one may estimate again the latter term as:

$$I_{\ell}(t) \ge j_1(t_0) + \operatorname{erfc}\left(\sqrt{\frac{3}{\mu t_0}}\right) \frac{\varepsilon}{4\sqrt{3}} \int_{t_0}^t \exp((b - s_{\infty})(t - t_0)) d\tilde{t}$$

= $j_1(t_0) + \operatorname{erfc}\left(\sqrt{\frac{3}{\mu t_0}}\right) \frac{\varepsilon}{4\sqrt{3}} \frac{(\exp((b - s_{\infty})(t - t_0)) - 1)}{(b - s_{\infty})},$

which leads to the inverse function:

$$t_{u} := \begin{cases} j_{1}^{-1}(\varrho_{0}) & \text{if } \varrho_{0} < j_{1}(t_{0}) \\ t_{0} + \frac{1}{(b-s_{\infty})} \ln \left(1 + \frac{4\sqrt{3}(b-s_{\infty})}{\varepsilon \operatorname{erfc}\left(\sqrt{\frac{3}{\mu t_{0}}}\right)} \left(\varrho_{0} - j_{1}(t_{0}) \right) \right) & \text{otherwise.} \end{cases}$$

when studying the limit of the latter expression when $\mu \to \infty$ one concludes that $j_1(t_0) \to 0$ and $\overline{t} := t_0 + \ln(1 + 4\sqrt{3}(b - s_\infty)\varrho_0/\varepsilon)$.

Corollary 1 Under hypotheses of Theorem 7, one might give another upper bound for t_{ρ_0} :

$$t_{\rho_0} \le t_u^{\infty} := \begin{cases} j_1^{-1}(\rho_0) & \text{if } \rho_0 < j_1(t_0), \\ j_2^{-1}(\rho_0) & \text{otherwise}, \end{cases}$$

where

$$t_0 = \frac{z}{2\sqrt{b - s_\infty}}, \quad j_1(t_0) := \frac{\varepsilon z}{4\mu e\sqrt{\pi}} \exp\left(-\frac{z\sqrt{b - s_\infty}}{\sqrt{\mu}}\right),$$

and

$$\begin{split} j_1^{-1}(\varrho_0) &:= \frac{z^2}{-2\ln\left(\frac{\varrho_0 4e\mu}{\varepsilon z}\right)},\\ j_2^{-1}(\varrho_0) &:= \left(\frac{(k+\frac{1}{2})k!\sqrt{\pi\mu}}{(b-s_\infty)^k}\exp((b-s_\infty)t_0)(\varrho_0-j_1(t_0)) + t_0^{k+\frac{1}{2}}\right)^{\frac{1}{k+\frac{1}{2}}}. \end{split}$$

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Fig. 4 The mutation rate μ is plotted in the log scale on the *x*-axis. The left (resp. *right*) figure displays the unbounded (resp. bounded) case. On the *left side*, one shows the upper and lower estimates of t_{ρ_0} , the time to reach ρ_0 , between these curves we computed numerically t_{I_l} and t_{I_u} in order to validate our estimates. On the *right*, one *plots* a numerical computation of $\mathcal{G}_{\varepsilon}$ solving (36) and the respective theoretical bounds

The proof simply takes into account that the Green function $\mathcal{G}_{b-s_{\infty}}$ in the bounded case is greater that the Green function in the case of the whole space, applying the same arguments as in Sect. 5.1.1 one concludes. These latter estimates shall improve the upper bound of t_{ρ_0} for μ small.

We plot in Fig. 4, right, the comparison between our bounds t_u , t_u^{∞} and t_l and the direct numerical simulations of $\mathcal{G}_{\varepsilon}$ solving (36) for a given set of data $(b, \varrho_0, \varepsilon, z)$ when μ varies.

5.2 Asymptotic expansion for large mutation rates

5.2.1 The formal result

When μ , the mutation rate dominates, one sets the decomposition :

$$\mathcal{N}_{\mu}(t,x) = \mathcal{N}_{0}(\mu t,x) + \frac{1}{\mu}\mathcal{N}_{1}(\mu t,x) + \cdots$$

and one gathers powers of $1/\mu$, separating the scales this leads to solve:

- At zero order one obtains the equation:

$$\partial_t \mathcal{N}_0 - \partial_{x^2}^2 \mathcal{N}_0 = 0$$

– While at *j*th order one writes:

$$\partial_t \mathcal{N}_j - \partial_{x^2}^2 \mathcal{N}_j = (b-s)\mathcal{N}_{j-1}, \quad j \in \mathbb{N}^*.$$

Using the spectral decomposition, one gets, for the zeroth order term, that it reads

$$\mathcal{N}_0(t, x) := \sum_{k \in \mathbb{N}} \overline{n}_I^k v_k(x) \exp(-\lambda_{0,k} t), \quad \overline{n}_I^k := < n_I, v_k >,$$

where the brackets denote the scalar product in $L^2((0, 1))$, and the eigenvectors are those of the homogeneous problem (13). This expression leads to exponentially decreasing modes and a mean that remains constant.

Denoting $\mathcal{N}_j := \sum_{k \in \mathbb{N}} \gamma_{j,k}(t) v_k(x)$ and $\Gamma_j(t) := (\gamma_{j,k}(t))_{k \in \mathbb{N}}$ one can write the modal equivalent of the equation above:

$$\dot{\Gamma}_j + \operatorname{diag}(\lambda)\Gamma_j = (b - \mathcal{M})\Gamma_{j-1}, \quad (\mathcal{M})_{ik} := \langle sv_i, v_k \rangle, \quad \Gamma_j(0) = 0, \ \forall j \ge 1.$$
(38)

where diag(λ) := diag($\lambda_{0,0}, \ldots, \lambda_{0,k}, \ldots$). As the dominant mode is the constant one, we prove below that

$$\gamma_{j,k} \sim \begin{cases} \overline{n}_I^0 \frac{(b - \mathcal{M}_{0,0})^{j} t^j}{j!} & \text{if } k = 0\\ 0 & \text{otherwise} \end{cases} + O(t^{j-1})$$

when *t* is large. This gives the asymptotic limit:

$$\mathcal{N}_{\mu} \sim \overline{n}_{I}^{0} \sum_{j} \frac{(b - \mathcal{M}_{0,0})^{j} t^{j}}{j! \mu^{j}} + R_{j} t^{j-1} = \overline{n}_{I}^{0} \exp\left(\frac{(b - \mathcal{M}_{0,0})t}{\mu}\right) + \dots$$

Returning to the original variables and after integration in time this gives that the final formula shall be

$$t_{\varrho_0} := \frac{1}{b - \mathcal{M}_{0,0}} \ln \left(1 + \frac{\rho_0 \left(b - \mathcal{M}_{0,0} \right)}{\left(\int_0^1 s(x) dx \right) \overline{n}_I^0} \right) + O\left(\frac{1}{\mu} \right).$$

Assuming that *s* fulfills hypotheses 3.2 this provides

$$t_{\varrho_0} := \frac{1}{b-\varepsilon} \ln\left(1 + \frac{\rho_0 \left(b-\varepsilon\right)}{\varepsilon \overline{n}_I^0}\right) + O\left(\frac{1}{\mu}\right). \tag{39}$$

5.2.2 Numerical simulations

We display in Fig. 5 for various values of μ , t_{ϱ_0} computed using direct numerical simulations with a random initial data n_I , compared with the value given by (39). In this particular case, the convergence occurs for values of μ larger then 1.

5.2.3 Mathematical proofs

Setting $\tilde{n}(t, x) := n(t, x) \exp(-(b - \mathcal{M}_{0,0})t)/\overline{n}_{I}^{0}$, with $\overline{n}_{I}^{0} := \int_{\Omega} n_{I}(x)dx$, it is solution of $\begin{cases} \partial_{t}\tilde{n} - \mu \partial_{x^{2}}^{2}\tilde{n} = (\mathcal{M}_{0,0} - s)\tilde{n} & (t, x) \in \mathcal{O}_{T}, \\ \partial_{x}\tilde{n}(t, x) = 0 & (t, x) \in \Sigma_{T}, \\ \tilde{n}(0, x) = n_{I}(x)/\overline{n}_{I}^{0} & (t, x) \in \{0\} \times \Omega \end{cases}$ (40)

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Fig. 5 Time to reach ρ_0 starting from a random initial condition, for various values of μ (in the logscale)

When μ , the mutation rate dominates one sets the decomposition:

$$\mathcal{N}_{\mu} = \mathcal{N}_0 + \frac{1}{\mu}\mathcal{N}_1 + \cdots$$

We suppose moreover that we focus on solutions for long times so that the time scaling should be

$$\mathcal{N}_{\mu}(\mu t, x) \sim \tilde{n}(t, x), \quad (t, x) \in \mathbb{R}_{+} \times (0, 1).$$

One then writes the asymptotic expansion of the equations with respect to $1/\mu$:

 $- \mathcal{O}(\mu)$ terms provide:

$$\begin{cases} \partial_t \mathcal{N}_0 - \partial_{x^2}^2 \mathcal{N}_0 = 0, & (t, x) \in \mathcal{O}_T, \\ \partial_x \mathcal{N}_0(t, x) = 0, & (t, x) \in \Sigma_T, \\ \mathcal{N}_0(0, x) = n_I(x)/\overline{n}_I^0, & (t, x) \in \{0\} \times \Omega. \end{cases}$$

– while $\mathcal{O}\left(\frac{1}{\mu^{j}}\right)$ terms give :

$$\begin{cases} \partial_t \mathcal{N}_j - \partial_{x^2}^2 \mathcal{N}_j = (\mathcal{M}_{0,0} - s) \mathcal{N}_{j-1}, & (t, x) \in \mathcal{O}_T, \\ \partial_x \mathcal{N}_j(t, x) = 0, & (t, x) \in \Sigma_T, \\ \mathcal{N}_j(0, x) = 0, & (t, x) \in \{0\} \times \Omega, \end{cases}$$

for $j \in \mathbb{N}^*$.

Proposition 3 For any $n_I \in L^2(\Omega)$, one has

$$-If j = 0$$

$$\|\mathcal{N}_0(t, \cdot) - 1\|_{L^2(0, 1)} \le c_0 \exp(-\pi^2 t), \quad \forall t \ge 0$$

- If $j \ge 1$ then

$$\|\mathcal{N}_{j}(t,\cdot)\|_{L^{2}(0,1)} \le c_{j}t^{j-1}, \quad \forall t \ge 0$$

where the constant c_i does not depend on t.

Proof The proof follows in two steps

- If j = 0 then the spectral decomposition gives directly the claim.
- For the second part we proceed by induction. If $j = 1, N_1$ solves

$$\partial_t \mathcal{N}_1 - \partial_{x^2}^2 \mathcal{N}_1 = (\mathcal{M}_{0,0} - s)\mathcal{N}_0$$

complemented with homogeneous Neumann boundary and zero initial conditions. Integrating in space shows that $\overline{\mathcal{N}}_1(t) := \int_{\Omega} \mathcal{N}_1(t, x) dx$ solves

$$\partial_t \overline{\mathcal{N}}_1 = \overline{(\mathcal{M}_{0,0} - s)\mathcal{N}_0} = \overline{(\mathcal{M}_{0,0} - s)\sum_{k \in \mathbb{N}} \gamma_{0,k} v_k} = -\sum_{k \neq 0} \mathcal{M}_{0,k} \gamma_{0,k} v_{k}$$

but the functions $\gamma_{0,k}$ are explicit and read: $\gamma_{0,k} = \overline{n}_I^k \exp(-\lambda_k t)$. Integrating in time and using Cauchy-Schwartz gives:

$$|\overline{\mathcal{N}}_{1}|^{2} \leq \left(\sup_{k \in \mathbb{N}} |\mathcal{M}_{0,k}|\right)^{2} \left|\sum_{k \neq 0} \overline{n}_{I}^{k} \left(\frac{(1 - \exp(-\lambda_{k}t))}{\lambda_{k}}\right)\right|^{2} \leq c \|s\|_{L^{\infty}(0,1)}^{2} \|n_{I}\|_{L^{2}(0,1)}^{2}$$

For the rest, we set $\underline{\mathcal{N}}_1 := \mathcal{N}_1 - \overline{\mathcal{N}}_1$, it solves

$$\partial_t \underline{\mathcal{N}}_1 - \partial_{x^2}^2 \underline{\mathcal{N}}_1 = (\mathcal{M}_{0,0} - s)\mathcal{N}_0 - \partial_t \overline{\mathcal{N}}_1,$$

which multiplied by \underline{N}_1 and integrated with respect to x reads :

$$\frac{1}{2}\partial_t \left\| \underline{\mathcal{N}}_1 \right\|_{L^2(0,1)}^2 + \lambda_1 \left\| \underline{\mathcal{N}}_1 \right\|_{L^2(0,1)}^2 \le c \left\| \mathcal{N}_0 \right\|_{L^2(0,1)} \left\| \underline{\mathcal{N}}_1 \right\|_{L^2(0,1)}^2$$

by Young's inequality and Gronwall's lemma, the second claim holds for j = 1.

– We suppose now that for $\ell \leq j - 1$ the property is true. Again for the constant mode one has

$$\partial_t \overline{\mathcal{N}}_j = \int_{\Omega} (\mathcal{M}_{0,0} - s) \mathcal{N}_{j-1} dx \le \|\mathcal{M}_{0,0} - s\|_{L^2(0,1)} \|\mathcal{N}_{j-1}\|_{L^2(0,1)} \le ct^{j-2},$$

where the last inequality comes from the induction hypothesis. Integrating the latter inequality in time provides the result for the zero mode. As above, one has

$$\partial_t \underline{\mathcal{N}}_j - \partial_{x^2}^2 \underline{\mathcal{N}}_j = (\mathcal{M}_{0,0} - s)\mathcal{N}_{j-1} - \partial_t \overline{\mathcal{N}}_j,$$

the latter term being constant in space and \underline{N}_j being of zero mean value provides after multiplication by \underline{N}_j and integration with respect to x that

$$\frac{1}{2}\partial_t \left\| \underline{\mathcal{N}}_j \right\|_{L^2(0,1)}^2 + \lambda_1 \left\| \underline{\mathcal{N}}_j \right\|_{L^2(0,1)}^2 \le 2 \|s\|_{L^\infty(0,1)} \left\| \mathcal{N}_{j-1} \right\|_{L^2(0,1)} \left\| \underline{\mathcal{N}}_j \right\|_{L^2(0,1)}$$

which again by Young inequality and Gronwall proves the second claim for $\ell = j$. This ends the induction argument.

5.2.4 Error estimates

We define the error $\mathcal{E}_N := n(t, x) \exp(-bt)/\overline{n}_I^0 - \exp(-\mathcal{M}_{0,0}t)\mathcal{N}_{\mu,N}$ where $\mathcal{N}_{\mu,N} := \sum_{j=0}^{j=N} \frac{1}{\mu^j} \mathcal{N}_j(\mu t, x)$, it solves

$$\begin{cases} \partial_t \mathcal{E}_N - \mu \partial_{x^2}^2 \mathcal{E}_N + s \mathcal{E}_N = \frac{(\mathcal{M}_{0,0} - s)}{\mu^N} \mathcal{N}_N(\mu t, x) \exp(-\mathcal{M}_{0,0} t), & (t, x) \in \mathcal{O}_T, \\ \partial_x \mathcal{E}_N = 0, & (t, x) \in \Sigma_T, \\ \mathcal{E}_N(0, x) = 0, & (t, x) \in \{0\} \times \Omega. \end{cases}$$

Theorem 9 If $n_I \in L^2(\Omega)$ and s satisfies hypotheses 3.2, one has for any fixed time *T* that

$$\|\mathcal{E}_N(T,\cdot)\|_{L^2(\Omega)} \le \frac{c}{\mu \mathcal{M}_{0,0}^N},$$

where the constant *c* is independent both on *T* and on μ .

Proof Multiplying the latter equation by \mathcal{E}_N and integrating with respect to x, one has

$$\partial_t \left\| \mathcal{E}_N \right\|_{L^2(\Omega)}^2 \le \frac{2c}{\mu^N} \left\| \mathcal{N}_N \right\|_{L^2(\Omega)} \left\| \mathcal{E}_N \right\|_{L^2(\Omega)} \exp(-\mathcal{M}_{0,0}t)$$

dividing both sides by $\sqrt{\|\mathcal{E}_N\|_{L^2(\Omega)}^2 + \delta}$, one gets:

$$\partial_t \sqrt{\|\mathcal{E}_N\|_{L^2(\Omega)}^2 + \delta} \le \frac{2c}{\mu^N} \frac{\|\mathcal{N}_N\|_{L^2(\Omega)} \|\mathcal{E}_N\|_{L^2(\Omega)}}{\sqrt{\|\mathcal{E}_N\|_{L^2(\Omega)}^2 + \delta}} \quad \exp(-\mathcal{M}_{0,0}t)$$
$$\le \frac{2c}{\mu^N} \|\mathcal{N}_N(\mu t, \cdot)\|_{L^2(\Omega)} \exp(-\mathcal{M}_{0,0}t)$$

here we use the estimates of Proposition 3 to conclude that

$$\sqrt{\left\|\mathcal{E}_{N}(t,\cdot)\right\|_{L^{2}(\Omega)}^{2}+\delta} \leq \sqrt{\delta} + \frac{c}{\mu} \int_{0}^{t} s^{N-1} \exp(-\mathcal{M}_{0,0}s) ds \leq \sqrt{\delta} + \frac{c}{\mu \mathcal{M}_{0,0}^{N}},$$

which finally gives the claim, since the result holds for any arbitrarily small $\delta > 0$.

Turning back to the original function n(t, x) solving (29), the consequence of results above is that

Corollary 2 $n_I \in L^2(\Omega)$ and s satisfies hypotheses 3.2, one has for any fixed time t that

$$\left\|n(t,\cdot)-\overline{n}_{I}^{0}\exp((b-\mathcal{M}_{0,0})t)\right\|_{L^{2}(\Omega)} \leq \frac{c|\overline{n}_{I}^{0}|}{\mu\mathcal{M}_{0,0}^{N}}\exp(bt).$$

where the constant c is independent both on t and on μ .

Remark 3 These results show that it is possible to compute the limit of *n* when μ is large. Nevertheless, the approximation is only first order accurate with respect to μ . An interesting and open question is whether one is able to construct a higher order approximation.

Theorem 10 If $n_I \in L^2(\Omega)$ and s satisfies hypotheses 3.2, for any $\eta > 0$, there exists a μ large enough, s.t. if

$$t_{\varrho_0} := \frac{1}{b - \mathcal{M}_{0,0}} \ln \left(1 + \frac{\rho_0 \left(b - \mathcal{M}_{0,0} \right)}{\left(\int_0^1 s(x) dx \right) \overline{n}_I^0} \right),$$

then $|\varrho_{\text{out}}(t_{\varrho_0}) - \varrho_0| \leq \eta$.

5.3 Asymptotic expansion for small mutation rates

On the other hand, if we consider μ small, one can decompose the solution of (29)

$$n(t,x) \sim \sum_{j \in \mathbb{N}} \mu^j \mathcal{N}_j(t,x)$$

where the different terms \mathcal{N}_i solve:

if
$$j = 0$$

$$\begin{cases}
\partial_t \mathcal{N}_0 = (b - s(x))\mathcal{N}_0, & (t, x) \in \mathcal{O}_T, \\
\mathcal{N}_0(0, x) = n_I(x).
\end{cases}$$
(41)

$$\begin{cases} \partial_t \mathcal{N}_j = (b - s(x))\mathcal{N}_j + \partial_{x^2}^2 \mathcal{N}_{j-1}, & (t, x) \in \mathcal{O}_T, \\ \mathcal{N}_j(0, x) = 0. \end{cases}$$
(42)

5.3.1 Formal computations

One can solve explicitly \mathcal{N}_0 , which reads:

$$\mathcal{N}_0(t, x) = \exp((b - s(x))t)n_I(x).$$

We compute a zero order approximation of ρ_{out} :

$$\varrho_{\text{out}} \sim \int_0^t \int_{\Omega} s(x) \exp((b - s(x))\tilde{t}) n_I(x) dx d\tilde{t} =: \rho_{\text{app}}(t).$$

For the particular case when s satisfies hypotheses 3.2, one recovers

$$\rho_{\rm app}(t) = \frac{\exp((b-1)t) - 1}{(b-1)} \int_{\Omega} s(x) n_I(x) dx,$$

which gives then an explicit formula for μ small

$$t_{\varrho_0} := \frac{1}{(b-1)} \log \left(1 + \frac{(b-1)\varrho_0}{\int_{\Omega} s(x) n_I(x) dx} \right).$$
(43)

5.3.2 Numerical simulations

We display, in Fig. 6, t_{ϱ_0} , the time to reach ϱ_0 , for various values of μ and for a given random initial data $n_I \in L^2(\Omega)$.

5.3.3 Rigorous proofs

Lemma 8 We suppose that s satisfies the hypotheses 3.2. For any given $\Psi \in L^2(\mathcal{O}_T)$ there exists a unique solution $\varphi \in C((0, T); H^1(\Omega)) \cap L^2((0, T); H^2(\Omega))$ solving

$$\begin{cases} \partial_t \varphi - \mu \partial_{x^2}^2 \varphi + s\varphi = \Psi(t, x), & (t, x) \in \mathcal{O}_T \\ \partial_x \varphi(t, x) = 0, & (t, x) \in \Sigma_T \\ \varphi(0, x) = 0, & (t, x) \in \{0\} \times \Omega, \end{cases}$$
(44)



Fig. 6 Time to reach ρ_0 starting from a random initial condition, for various values of μ , we use again the same code as in Fig. 5

and one has continuity with respect to the $L^2((0,T); H^2(\Omega))$ norm:

$$\|\varphi\|_{L^2((0,T);H^2(\Omega))} \le C \|\Psi\|_{L^2(\mathcal{O}_T)}.$$

The proof, based on the Galerkin decomposition, is classical and can be found in Ladyženskaja et al. (1968) Chap III, Sect. 6, p. 172–178.

Theorem 11 If $n_I \in L^2(\Omega)$ and s satisfies hypotheses 3.2, one has the $L^2(\mathcal{O}_T)$ -error estimates:

$$\|n - \mathcal{N}_0\|_{L^2(\mathcal{O}_T)} \le \mu \sqrt{\frac{\exp(2bT) - 1}{2b}} \|n_I\|_{L^2(\Omega)}$$

Proof We rescale the problem (29) and the asymptotic expansion so to drop the damping term *b*, in the respective equations. Let define $Y := C^0((0, T); H^1(\Omega)) \cap L^2((0, T); H^2(\Omega))$. Considering (41), the equation is satisfied also in the weak sense namely, for any $\varphi \in Y$, one has

$$\left[\int_{\Omega} \mathcal{N}_0(t,x)\varphi(t,x)dx\right]_{t=0}^{t=T} - \int_{\mathcal{O}_T} \mathcal{N}_0(t,x)\left(\partial_t \varphi - s(x)\varphi\right)dxdt = 0.$$

On the other hand, one has also that *n* solving (29), satisfies as well for any $\varphi \in Y$,

$$\left[\int_{\Omega} n(t,x)\varphi(t,x)dx\right]_{t=0}^{t=T} - \int_{\mathcal{O}_T} n(t,x)\left(\partial_t \varphi + \mu \partial_{x^2}^2 \varphi - s(x)\varphi\right)dxdt = 0.$$

This gives when setting $\mathcal{E}(t, x) := n(t, x) - \mathcal{N}_0(t, x)$, for any $\varphi \in Y$

$$\begin{split} \left[\int_{\Omega} \mathcal{E}(t,x)\varphi(t,x)dx \right]_{t=0}^{t=T} &- \int_{\mathcal{O}_{T}} \mathcal{E}(t,x) \left(\partial_{t}\varphi + \mu \partial_{x^{2}}^{2}\varphi - s(x)\varphi \right) dxdt \\ &= \mu \int_{\mathcal{O}_{T}} \mathcal{N}_{0}(t,x) \partial_{x^{2}}^{2}\varphi(t,x) dxdt. \end{split}$$

For any $\Psi \in L^2(\mathcal{O}_T)$, there exists $\varphi \in Y$, the forward form solving (44). Now let insert the backward expression of φ in the latter weak form, this gives:

$$\int_{\mathcal{O}_T} \mathcal{E}(t,x)\Psi(t,x)dxdt \leq \mu \|\mathcal{N}_0\|_{L^2(\mathcal{O}_T)} \left\|\partial_{x^2}^2\varphi\right\|_{L^2(\mathcal{O}_T)} \leq C\mu \|\mathcal{N}_0\|_{L^2(\mathcal{O}_T)} \|\Psi\|_{L^2(\mathcal{O}_T)},$$

which holds for any $\Psi \in L^2(\mathcal{O}_T)$. Taking the supremum over all functions in this latter space provides the bound $\|\mathcal{E}\|_{L^2(\mathcal{O}_T)} \leq C\mu \|\mathcal{N}_0\|_{L^2(\mathcal{O}_T)}$. Turning back to the correct scaling with respect to the damping term *b*, one recovers the claim.

Corollary 3 Under the same hypotheses as in Theorem (11), the time t_{ϱ_0} being defined in (43), one has the error estimates : there exists a constant $C(s, n_I, b) > 0$, independent of μ , s.t.

 $|\varrho_{\text{out}}(t_{\varrho_0}) - \varrho_0| \le \mu C(s, n_I, b) ||n_I||_{L^2(\Omega)} ||s||_{L^2(\Omega)}.$

The proof comes easily combining a triangular inequality as in the proof of Theorem 4, results above and the definition of t_{ϱ_0} . As the latter time does not depend on μ , the claim follows straightforwardly.

References

- Abbas AK, Lichtman AH, Pillai S (2012) Basic immunology: functions and disorders of the immune system. Elsevier Health Sciences, USA
- Balelli I, Milisic V, Wainrib G, Branching random walks on binary strings and application to adaptive immunity preprint, submitted
- Balelli I, Milisic V, Wainrib G, Random walks on binary strings applied to the somatic hypermutation of B-cells preprint, submitted
- Bendali A, Huard A, Tizaoui A, Tordeux S, Vila J-P (2009) Asymptotic expansions of the eigenvalues of a 2-D boundary-value problem relative to two cavities linked by a hole of small size. CR Math Acad Sci Paris 347(19–20):1147–1152. doi:10.1016/j.crma.2009.09.005
- Bendali A, Tizaoui A, Tordeux S, Vila JP (2008) Second-order asymptotic expansion for an eigenvalue set in domain with small iris, in: Progress in industrial mathematics at ECMI, Vol. 15 of Math. Ind., Springer, Heidelberg, 2010, pp 715–720. doi:10.1007/978-3-642-12110-4_115
- Bürger R (2000) The mathematical theory of selection, recombination, and mutation, vol 228. Wiley, Chichester
- Coddington EA, Levinson N (1955) Theory of ordinary differential equations. McGraw-Hill Book Company Inc, New York
- Dautray R, Lions J-L (1990) Mathematical analysis and numerical methods for science and technology, vol 3. Springer-Verlag, Berlin
- Diaz M, Casali P (2002) Somatic immunoglobulin hypermutation. Curr opin Immunol 14(2):235-240
- Dunn-Walters DK, Belelovsky A, Edelman H, Banerjee M, Mehr R (2002) The dynamics of germinal centre selection as measured by graph-theoretical analysis of mutational lineage trees. Clin Dev Immunol 9(4):233–243
- Faro J, Or-Guil M (2013) How oligoclonal are germinal centers? A new method for estimating clonal diversity from immunohistological sections. BMC Bioinform 14(6):S8. doi:10.1186/1471-2105-14-S6-S8
- Fisher RA (1999) The genetical theory of natural selection, variorum Edition, Oxford University Press, Oxford (revised reprint of the 1930 original, Edited, with a foreword and notes, by J. H. Bennett)
- Gitlin AD, Shulman Z, Nussenzweig MC (2014) Clonal selection in the germinal centre by regulated proliferation and hypermutation. Nature 509:637–640
- Kelsoe G (1996) Life and death in germinal centers (redux). Immunity 4(2):107-111
- Kepler TB, Perelson AS (1993) Cyclic re-entry of germinal center b cells and the efficiency of affinity maturation. Immunol Today 14(8):412–415
- Ladyženskaja OA, Solonnikov VA, Ural'ceva NN (1968) Linear and quasilinear equations of parabolic type, Translated from the Russian by S. Smith. Translations of Mathematical Monographs, Vol. 23, American Mathematical Society, Providence, R.I
- Meyer-Hermann M, Deutsch A, Or-Guil M (2001) Recycling probability and dynamical properties of germinal center reactions. J Theor Biol 210(3):265–285
- Meyer-Hermann M (2002) A mathematical model for the germinal center morphology and affinity maturation. J Theor Biol 216(3):273–300
- Neuberger MS, Ehrenstein MK, Klix N, Jolly CJ, Yélamos J, Rada C, Milstein C (1998) Monitoring and interpreting the intrinsic features of somatic hypermutation. Immunol Rev 162(1):107–116
- Nowak MA (2006) Evolutionary dynamics. Harvard University Press, Massachusetts

Perelson AS, Weisbuch G (1997) Immunology for physicists. Rev Mod Phys 69:1219–1268. doi:10.1103/ RevModPhys.69.1219

Perthame B (2006) Transport equations in biology. Springer, Berlin

Swerdlin N, Cohen IR, Harel D (2008) The lymph node b cell immune response: dynamic analysis in-silico. Proc IEEE 96(8):1421–1443

Tarlinton DM (2014) Immunology: to affinity and beyond. Nature 509(7502):573-574

Teng G, Papavasiliou FN (2007) Immunoglobulin somatic hypermutation. Annu Rev Genet 41:107–120 Victora GD (2014) Snapshot: the germinal center reaction. Cell 159(3):700–700

- Wabl M, Cascalho M, Steinberg C (1999) Hypermutation in antibody affinity maturation. Current opinion in immunology 11(2):186–189
- Zettl A (2005) Sturm-Liouville theory, vol 121 of mathematical surveys and monographs. American Mathematical Society, Providence, RI