

Origin of the spontaneous oscillations in a simplified coagulation-fragmentation system driven by a source

Jean-Yves Fortin

▶ To cite this version:

Jean-Yves Fortin. Origin of the spontaneous oscillations in a simplified coagulation-fragmentation system driven by a source. Journal of Physics A: Mathematical and Theoretical, 2022, 55, pp.485003. 10.1088/1751-8121/aca572. hal-03709590

HAL Id: hal-03709590

https://hal.science/hal-03709590

Submitted on 30 Jun 2022

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers. L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

Origin of the spontaneous oscillations in a simplified coagulation-fragmentation system driven by a source

Jean-Yves Fortin

Laboratoire de Physique et Chimie Théoriques, CNRS UMR 7019, Université de Lorraine, F-54000 Nancy, France

E-mail: jean-yves.fortin@univ-lorraine.fr

Abstract. We consider a system of aggregated clusters of particles, subjected to coagulation and fragmentation processes with mass dependent rates. Each monomer particle can aggregate with larger clusters, and each cluster can fragment into individual monomers with a rate directly proportional to the aggregation rate. The dynamics of the cluster densities is governed by a set of Smoluchowski equations, and we consider the addition of a source of monomers at constant rate. The whole dynamics can be reduced to solving a unique non-linear differential equation which displays self-oscillations in a specific range of parameters, and for a number of distinct clusters in the system large enough. This collective phenomenon is due to the presence of a fluctuating damping coefficient and is closely related to the Liénard self-oscillation mechanism observed in a more general class of physical systems such as the van der Pol oscillator.

PACS numbers: 05.20.Dd,05.45.-a,36.40.Qv,36.40.Sx

1. Introduction

The processes of aggregation and fragmentation are found in a large variety of systems at different levels of scale. We can mention various physical phenomena that fit into this class of problems, such as foam dynamics with coalescence of soap bubbles, the stability of atmospheric particles of different sizes which aggregate due to the van der Waals forces, or even internet large scale networks where nodes can acquire additional links according to their attachment preference weight [1]. On larger scales, it is known that aggregation plays an important role in the formation of planetary rings [2] or dust clouds induced by gravitational forces, and in the formation of galaxy clusters [3, 4]. In contrast, fragmentation acts as a counterweight to the aggregation phenomenon, leading in many cases to a stationary state for the mass or size distribution of the particle clusters. The mechanisms of fragmentation have diverse origins. For example, the shattering of heavy particles that collide leads to smaller particles, or tidal forces between large masses can break them apart because of the inhomogeneity of the gravitation between different parts of the same body. Thermal fluctuations are also responsible for the desegregation of large polymers in solutions.

A general kinetic theory that incorporates both the aggregation and fragmentation is usually constructed from a set of Smoluchowski type equations that govern the dynamics of the mass densities of each individual object. In the long time limit, the steady state distribution exhibits in many models, according to the type of kernel chosen for the aggregation or fragmentation processes, a power law behavior associated with an exponential decay [5, 6, 7, 8, 9]. A summary of several analytical solutions for the density distribution is presented in Table 1, for different kinds of kernels and eventually in the presence of a particle source. Apart from these time independent steady states, leading to an equilibrium distribution of masses, one can also observe collective oscillations in a narrow window of parameter range [10, 11, 12], see Table 2. However it is not yet clear how these oscillations arise from the complex dynamics involving different clusters of particles, since it is a highly non-linear process and the Smoluchowski equations are usually quadratic in the densities and exact solutions can not be found in general, except for small sets of equations [13, 14], and this is an ongoing challenge in the field of mathematics.

In this paper we would like to investigate analytically the nature and origin of these oscillations in a simplified model of aggregation and fragmentation process. Here we consider a set of $N \gg 1$ possible distinct masses $M_k = kM_1$, $k \leq N$, with their density n_k , and we restrict in this manuscript the dynamics to the aggregation process between the monomer particles and bigger masses $M_1 + M_k \to M_{k+1}$, and the fragmentation of large masses into the smallest units $M_1 + M_k \to (k+1)M_1$, see figure 1. We also add a source of monomers M_1 at constant rate and eventually a sink for the largest clusters M_N , which conserves the total mass of the system.

The manuscript, which tries to detail the dynamics of this model, is organized as follow. In section 2, we describe the general model and write the set of Smoluchowski master equations for the dynamics of each cluster density. These equations can be reduced to a set of self-consistent integral equations, depending solely on the solution for the monomer density which satisfies a non-linear integral equation of Volterra type. The structure of these integrals is dependent on the eigenstates of the main coefficient matrix with eigenvalues in the complex plane. As a first simple example, we demonstrate in section 3 that the monomer density can only be positive and develops a complex time signal. In section 4, we consider the approximation case of a single harmonic oscillator with negative damping for representing the eigenstate whose eigenvalue has the largest positive real part and we rewrite the Volterra equation for the monomer density as a second order non-linear ordinary differential equation (ODE). The monomer density reaches either a constant value in the long time limit or displays steady oscillations, depending on the input parameters of the model. The frequency of these oscillations is independent on the current and can be computed perturbatively around the oscillations onset. In section 5, we extend this method to the main model in order to obtain a more precise value of the oscillation frequency, and in the conclusion we explain why this oscillating phenomenon is related to the Liénard self-oscillations observed in non-linear mechanical systems with a negative damping.

Table 1: Asymptotic distributions for several models of aggregation-fragmentation in the steady state. The notations are the following: Aggregation kernel K_{ij} , fragmentation kernel F_{ij} , spontaneous fragmentation rate F_k , and source of k-mass clusters I_k . $\theta = \nu - \mu \ge 0$ and $\beta = \mu + \nu$. The mass density is normalized to unity.

Kernels K_{ij} , F_{ij} , F_k , and current I_k	Asymptotic distribution	References
$K_{ij} = 1$, constant kernel, $F_{ij} = 0$	$n_k(t) \simeq 4t^{-2}e^{-2k/t}$	[15]
$K_{ij} = (i+j)/2$, additive kernel, $F_{ij} = 0$	$n_k(t) \simeq e^{-2t} g(ke^{-t}), g(z) = e^{-z/2} (2\pi z^3)^{-1/2}$	[15]
$K_{ij} = ij$, multiplicative kernel, $F_{ij} = 0$	$n_k(t) \simeq (2\pi)^{-1/2} k^{-5/2} t^{-1}$	[15]
$K_{ij}=ij, F_{ij}=0, F_k=\nu k, \nu \geq 2 \text{ (gelation for } \nu \leq 2\text{)}$	$n_k \simeq \frac{(\nu-2)k + (1+\alpha)}{2\sqrt{\pi}} e^{-k\log(\nu/2)} k^{-7/2}$	[16]
$K_{ij} = i^{\mu}j^{\nu} + i^{\nu}j^{\mu}, F_{ij} = 0, I_k = I\delta_{k1}, \theta < 1 \text{ and } \beta < 1$	$n_k \simeq \left(\frac{I(1-\theta^2)\cos(\pi\theta/2)}{4\pi}\right)^{1/2} k^{-(3+\beta)/2}$	[17, 18]
$K_{ij} = i^{\mu}j^{\nu} + i^{\nu}j^{\mu}, \ F_{ij} = \lambda K_{ij}, \ \mu = \nu = \beta/2$	$n_k \simeq n_1 \frac{(1+\lambda)^2}{\sqrt{\pi}(1+2\lambda)} e^{-k} k^{-(3+\beta)/2}$	[11]
$K_{ij} = 1, F_{ij} = \lambda, F_k = \nu$	$n_k \simeq c(4\pi)^{-1/2}e^{-\nu k/2}k^{-3/2},$	
	$c = (\lambda - \nu + \eta)/(1 + 2\lambda)$	
	$\eta = \sqrt{(\lambda - \nu)^2 + 2\nu(1 + 2\lambda)}$	[9]
$K_{ij} = (ij)^{\mu}, F_{ij} = \lambda, F_k = 0$	$n_k \simeq \frac{\lambda^{1-2\mu}}{\Gamma(1/2-\mu)} e^{-\lambda^2 k} k^{-3/2-\mu}$	[9]
$K_{1i} = K_{i1} = 1, F_{1i} = F_{i1} = \lambda, \lambda > \lambda_c$	$n_k \simeq \frac{\lambda}{1+\lambda} e^{-k\log(1+\lambda)}$	[8]
$K_{1i} = K_{i1} = i, F_{1i} = F_{i1} = \lambda i, \lambda > \lambda_c$	$n_k \simeq \lambda e^{-k\log(1+\lambda)} k^{-1}$	[8]

Table 2: Models with stable oscillations

Kernels K_{ij} , F_{ij} , and current I_k	Parameter range for stable oscillations	References
$K_{ij} = i^{\mu}j^{\nu} + i^{\nu}j^{\mu}, I_k = I\delta_{k1}, \text{ mass cutoff } M \gg 1$	$ \theta = \nu - \mu > 1$	[10, 11]
$K_{ij} = i^{-a}j^a + i^aj^{-a}, F_{ij} = \lambda K_{ij}, I_k = 0$	Brownian kernel $a = \nu = -\mu, a \simeq 1^-, \lambda \ll 1$	[12]
$K_{ij} = i^{\mu}j^{\nu} + i^{\nu}j^{\mu}, F_{ij} = \lambda K_{ij}, I_k = 0$	$1 < \theta = \nu - \mu \simeq 2^- \ (1/2 < a \simeq 1^-)$	
	and $\lambda < \lambda_c(\theta) \ll 1$	[12]

2. Main model

We consider a set of clusters containing k particles, $1 \le k \le N$, each with a density n_k . Only one-particle clusters react with other clusters of size $k \ge 1$, see figure 1. A particle can aggregate with a cluster of size k to form a cluster of size k + 1 with probability K_k . Or it can shatter this cluster entirely into k single particles with probability λK_k , as considered in previous works [12]. Usually it is assumed that $\lambda \ll 1$. The rates K_k grow according to the cluster size k, as a power law with exponent a, such that $K_k \propto k^a$. A source of monomer particles is finally introduced with

rate μ , and a sink term with rate $-\mu/N$ for the N-particle clusters in order to conserve the total mass. The Smoluchowski equations therefore read

$$\frac{\partial n_1}{\partial t} = -2n_1^2 - \sum_{k=2}^{N-1} (1 - k\lambda) K_k n_1 n_k + N\lambda K_N n_1 n_N + \mu,$$

$$\frac{\partial n_k}{\partial t} = K_{k-1} n_1 n_{k-1} - (1 + \lambda) K_k n_1 n_k, \ 2 \le k \le N - 1$$

$$\frac{\partial n_N}{\partial t} = K_{N-1} n_1 n_{N-1} - \lambda K_N n_1 n_N - \mu/N.$$
(1)

This set of equations satisfies the balance relation for the mass current $\partial_t \sum_{k\geq 1} kn_k = 0$. It is convenient to redefine the cluster densities by introducing the rescaled variables $v_k = K_k n_k$ and a new definition for the time $\tau = \tau(t)$ such as $d\tau = n_1 dt$. We obtain in this case a series of quasi-linear ODEs

$$\frac{\partial v_1}{\partial \tau} = -2v_1 - \sum_{k=2}^{N-1} (1 - k\lambda)v_k + N\lambda v_N + \mu/v_1,$$

$$K_k^{-1} \frac{\partial v_k}{\partial \tau} = v_{k-1} - (1 + \lambda)v_k, \ 2 \le k \le N - 1$$

$$K_N^{-1} \frac{\partial v_N}{\partial \tau} = v_{N-1} - \lambda v_N - \mu(Nv_1)^{-1}.$$
(2)

2.1. Algebraic analysis

In order to solve formally the system (2), we rewrite it in the form $\partial_{\tau} \mathbf{v} = A\mathbf{v} + \mathbf{j}$, with $\mathbf{v} = (v_1, \dots, v_N)$, A the $N \times N$ coefficient matrix of the system, and $\mathbf{j} = \mu(1/v_1, 0, \dots, 0, -K_N(Nv_1)^{-1})$

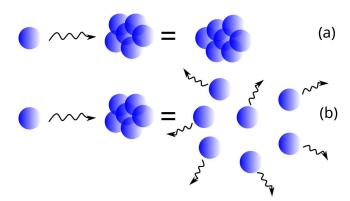


Figure 1: Aggregation and fragmentation process for the main model. A one-particle collides and aggregates with a k-particle cluster of size k with probability K_k , or fragment it into k individual particles with probability λK_k .

the current vector. The matrix A is given explicitly by

$$A = \begin{pmatrix} -2 & -1 + 2\lambda & -1 + 3\lambda & \cdots & -1 + (N-1)\lambda & N\lambda \\ K_2 & -(1+\lambda)K_2 & 0 & 0 & \cdots & 0 \\ 0 & K_3 & -(1+\lambda)K_3 & 0 & \cdots & 0 \\ \vdots & 0 & K_4 & -(1+\lambda)K_4 & 0 & 0 \\ 0 & 0 & 0 & \ddots & \ddots & 0 & 0 \\ 0 & \cdots & 0 & K_{N-1} & -(1+\lambda)K_{N-1} & 0 \\ 0 & 0 & \cdots & 0 & K_N & -\lambda K_N \end{pmatrix}.$$
(3)

The general ODE for $\boldsymbol{v}(\tau)$ can be solved formally using the diagonalization of A. This leads directly to the integral solution

$$\boldsymbol{v}(\tau) = \int_0^{\tau} P e^{D(\tau - \tau')} P^{-1} \boldsymbol{j}(\tau') d\tau' + P e^{D\tau} P^{-1} \boldsymbol{v}(0), \tag{4}$$

where $A = PDP^{-1}$, D is the diagonal matrix containing the complex eigenvalues $D = \text{diag}(z_1, \dots, z_N)$, and P the matrix of the eigenvectors. As initial condition, we take a system consisting of only monomer particles, $\mathbf{v}(0) = (v_1(0), 0, \dots, 0)$, with $v_1(0) = n_1(0) = M_0$ and M_0 being the total initial mass that we can take to be unity for simplification. The eigenvalues z_i are roots of the polynomial equation

$$(z+2)\prod_{k=2}^{N-1}(1+\lambda+zK_k^{-1})+\sum_{k=2}^{N-1}(1-k\lambda)\prod_{l=k+1}^{N-1}(1+\lambda+zK_l^{-1})=\frac{\lambda N}{\lambda+zK_N^{-1}},$$
 (5)

with the convention that when k = N - 1 in the sum, the product inside the sum is equal to one. We can check that the eigenvalue z = 0 is always solution of this polynomial, and corresponds physically to a steady state. In the limit $N \gg 1$, with $K_k = k^a$, $\lambda > 0$ and a > 0, the product $\prod_{k=2}^{N-1} (1+\lambda+zK_k^{-1})$ is generally diverging and the secular equation (5) tends to the implicit series

$$1 = \sum_{k=1}^{\infty} (\lambda k - 1) \prod_{l=1}^{k} (1 + \lambda + z K_l^{-1})^{-1}.$$
 (6)

This equation has real solutions near the poles $z = -(1 + \lambda)K_k < 0$, as well as the steady state solution z = 0. However it is difficult to extract complex solutions in this limit, even numerically. For the simplest case when N is finite and a = 0, or $K_k = 1$, the equation (5) has N - 1 solutions located on the unit circle of center $-1 - \lambda$, with the value $z = -\lambda$ excluded, in addition to the trivial solution z = 0

$$z_k = -1 - \lambda + e^{2i\pi k/N}, \ k = 1, \dots, N - 1, \text{ and } z_0 = 0.$$
 (7)

All these solutions except the zero solution have strictly negative real part. To obtain complex solutions with a strictly positive part, a has to be positive, and N sufficiently large. In figure 2 are represented the complex eigenvalues for different sizes N. The set of solutions deviates from the circle (7), and one observes that two complex conjugated eigenvalues converge to a finite value with a positive part in the asymptotic limit, see the inset of figure 2. In particular, this occurs

when N > 10. When N = 10, all the non-zero eigenvalues have negative real parts, and in this case we do not observe any oscillatory behavior, as we will see in the following. A necessary condition, but a priori not sufficient, to observe oscillations is to have at least one pair of eigenvalues with a strictly positive real part, which occurs already for N = 20 in figure 2. A more precise analysis will be made in section 5. Each component of the vector \boldsymbol{v} solution of the system (4), and satisfying the previous initial condition, can be expressed as

$$v_k(\tau) = \mu \int_0^\tau \frac{Q_k(\tau - \tau')}{v_1(\tau')} d\tau' + \tilde{Q}_k(\tau), \tag{8}$$

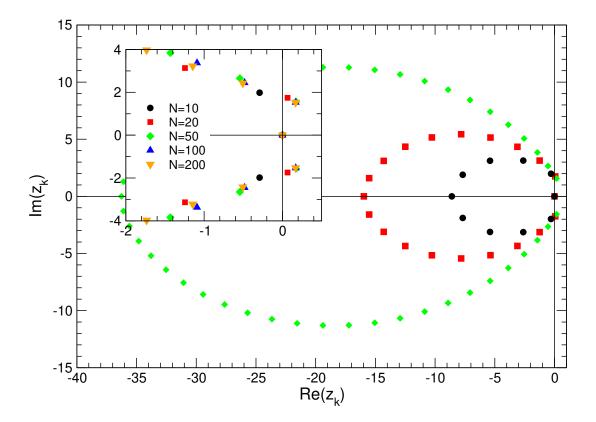


Figure 2: Set of eigenvalues of the matrix A, with N = 10, 20, 50, 100, 200, exponent a = 0.9, and fragmentation factor $\lambda = 0.04$. Each eigenvalue comes in conjugate pair in the complex plane, and the zero eigenvalue is always solution of the polynomial (5). In inset are represented the eigenvalues close to the origin, which converge as N becomes large. Only two of them have a strictly positive real part.

with the kernel functions Q_k and \tilde{Q}_k defined by the eigenvectors of the matrix A

$$Q_k(\tau) = \sum_{l=1}^{N} P_{kl} \left((P^{-1})_{l1} - \frac{K_N}{N} (P^{-1})_{lN} \right) e^{z_l \tau},$$

$$\tilde{Q}_k(\tau) = \sum_{l=1}^{N} P_{kl} (P^{-1})_{l1} e^{z_l \tau},$$
(9)

and initial conditions $Q_k(0) = \delta_{k,1} - (K_N/N)\delta_{k,N}$, $\tilde{Q}_k(0) = \delta_{k,1}$. The behavior of $Q_k(\tau)$ or $\tilde{Q}_k(\tau)$ is dominated after some transit time by the set of conjugate eigenvalues with the largest real part. If this one has strictly positive real and non zero imaginary parts, Q_k and \tilde{Q}_k will oscillate and diverge exponentially. Even if the Q_k s and \tilde{Q}_k s are exponentially diverging with time, we will show that the density solutions are finite. The above equations are all determined by $v_1(\tau)$, which satisfies the following self-consistent functional equation

$$v_1(\tau) = \mu \int_0^{\tau} \frac{Q_1(\tau - \tau')}{v_1(\tau')} d\tau' + \tilde{Q}_1(\tau), \tag{10}$$

and which is usually called a Volterra integral equation of the second kind in the Hammerstein form [19]. It is in general difficult to obtain an exact solution for any arbitrary function Q_1 as the equation (10) is explicitly a non-linear form and generally non-local. The Laplace transform can be applied to (10) since it involves a convolution product, however this does not give any more information as there is no useful relations between the Laplace transforms of the function $v_1(\tau)$ and its reciprocal $v_1^{-1}(\tau)$.

2.2. Steady state solution

Before analyzing the time evolution of the equation $\partial_{\tau} \mathbf{v} = A\mathbf{v} + \mathbf{j}$, we can search for constant solutions in the long time limit. Since A is singular, solving directly (2) with $\partial_{\tau} v_i = 0$ does not lead to any definite answer. We consider instead the solution of the equation $A\mathbf{v} = 0$, which is a one-dimensional vector, and after some computation we find that the kernel $\ker(A)$ is generated by the vector

$$\mathbf{v} = v_1 \left(1, (1+\lambda)^{-1}, (1+\lambda)^{-2}, \cdots, (1+\lambda)^{2-N}, \lambda^{-1} (1+\lambda)^{2-N} \right)^{\top}, \tag{11}$$

where v_1 is any arbitrary constant. It is easy to show that if we add in the last component of this vector the term $-\mu(N\lambda v_1)^{-1}$, then the new vector satisfies $A\mathbf{v} + \mathbf{j} = 0$. In order to determine v_1 , we use the mass conservation equation $\sum_{k=1}^{N} kv_k/K_k = 1$ and obtain a second order equation for v_1

$$\sum_{k=1}^{N-1} \frac{kv_1}{K_k(1+\lambda)^{k-1}} + \frac{Nv_1}{K_N\lambda(1+\lambda)^{N-2}} - \frac{\mu}{K_N\lambda v_1} = 1.$$
 (12)

This equation always possesses only one positive solution v_1^* , and in particular, for $N \gg 1$, we obtain that $v_1^{*-1} \simeq \sum_{k\geq 1} k^{1-a}/(1+\lambda)^{k-1}$, with the notation $K_k = k^a$. This solution is independent

of μ in this limit. The densities n_k then satisfy the equilibrium values

$$n_k = \frac{k^{-a}}{(1+\lambda)^{k-1}} v_1^* \simeq \frac{e^{-\lambda(k-1)}}{k^a} v_1^*. \tag{13}$$

2.3. Numerical analysis

Numerically, we can solve the system of equations (2) with a 4^{th} order Runge-Kutta algorithm for N different densities v_k as a function of the effective time τ . In the following we consider only the time τ since the real time t is an increasing and monotonic function of τ . Indeed n_1 is found to be always strictly positive. In figure 3 is represented the first 10 densities $v_{k=1,\dots,10}$ when N=100 for the set of parameters a=0.9, $\lambda=0.04$, and $\mu=0.002$. They all display collective oscillations, with a retarded phase shift between each consecutive curve. For this number N, the largest pair of eigenvalues of the matrix A is found to be approximately equal to $z_{max}^{\pm} \simeq 0.169 \pm 1.531i$, and has therefore a strictly positive real part. All the solutions of system (2) are oscillating at a frequency

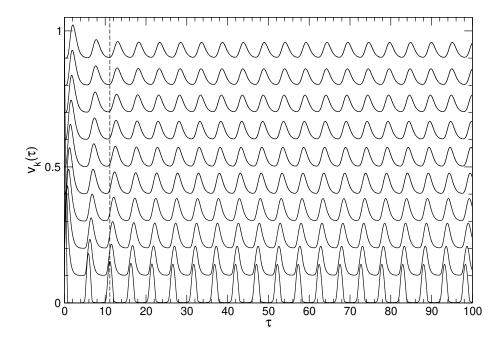


Figure 3: Numerical solution of the system (2) for the parameters N = 100, a = 0.9, $\lambda = 0.04$, and $\mu = 2.10^{-3}$. The successive oscillating densities from $v_1(\tau)$ (bottom) to $v_{10}(\tau)$ (top) are shown and shifted upwards for more clarity. The vertical dashed line shows the progressive phase shift between the densities.

 Ω estimated after Fourier transform of the v_1 signal to be equal to $\Omega \simeq 1.226$, which is lower than the imaginary part $\text{Im}(z_{max}^+)$. We therefore can not relate directly the observed frequency to $\text{Im}(z_{max}^+)$, which is a natural frequency of the problem, and a more detailed analysis has to be performed in the next sections.

3. Simplified model

In order to investigate the previous ODE (10), and show in particular that v_1 is always a positive quantity with time, we consider a very simplified kernel version of $Q_1(\tau)$ and $\tilde{Q}_1(\tau)$ which is represented by the general sawtooth or triangular function with equidistant time steps $\tau_{n\geq 0}$ such that $T = \tau_{n+1} - \tau_n$ with $\tau_0 = 0$, and increasing slope coefficients $\alpha_n > 0$

$$Q_1(\tau) = 1 - \alpha_0 \tau, \ 0 \le \tau \le T$$

$$Q_1(\tau) = Q_1(\tau_n) + (-1)^{n+1} \alpha_n (\tau - \tau_n), \ \tau_n \le \tau \le \tau_{n+1},$$
(14)

where the extrema are given by $Q_1(\tau_n) = 1 + T \sum_{k=0}^{n-1} (-1)^{k+1} \alpha_k$. Replacing Q_1 and \tilde{Q}_1 in equation (10) by (14), and differentiating with respect to τ , we obtain the following first order and non-local integro-differential equations for each time interval

$$v_1'(\tau) = \frac{\mu}{v_1(\tau)} - \mu \alpha_0 \int_{\tau - \tau_1}^{\tau} \frac{d\tau'}{v_1(\tau')} + \mu \alpha_1 \int_{\tau - \tau_2}^{\tau - \tau_1} \frac{d\tau'}{v_1(\tau')} - \cdots + (-1)^{n+1} \mu \alpha_n \int_0^{\tau - \tau_n} \frac{d\tau'}{v_1(\tau')} + (-1)^{n+1} \alpha_n, \ \tau_n \le \tau \le \tau_{n+1}.$$
(15)

Deriving a second time with respect to τ , we obtain a set of piecewise second order ODEs

$$v_1''(\tau) = -\frac{\mu v_1'(\tau)}{v_1^2(\tau)} - \frac{\mu \alpha_0}{v_1(\tau)} + \frac{\mu(\alpha_0 + \alpha_1)}{v_1(\tau - \tau_1)} - \frac{\mu(\alpha_1 + \alpha_2)}{v_1(\tau - \tau_2)} + \cdots + (-1)^{n+1} \frac{\mu(\alpha_{n-1} + \alpha_n)}{v_1(\tau - \tau_n)}, \ \tau_n \le \tau \le \tau_{n+1}.$$

$$(16)$$

Each solution on the intervals $\tau_n \leq \tau \leq \tau_{n+1}$ depends on the previous solutions $v_1(\tau - \tau_k)$ for $k = 0, \dots, n$ and can be recursively computed from the first solution in the interval $0 \leq \tau \leq T$ which satisfies

$$v_1^2(\tau)v_1''(\tau) = -\mu \left(v_1'(\tau) + \alpha_0 v_1(\tau)\right). \tag{17}$$

Initial conditions are $v_1(0) = 1$ and $v_1'(0) = \mu - \alpha_0 \ddagger$. Solving this system of ODEs requires the continuity of v_1 at the interval boundaries: $v_1(\tau_1^+) = v_1(\tau_1^-)$. There also exists a relation between the first derivatives at the boundaries $v_1'(\tau_1^+) = v_1'(\tau_1^-) + \alpha_0 + \alpha_1$, which can be derived directly from equation (15). We use the following generalized conditions at $\tau = \tau_n$:

$$v_1(\tau_n^+) = v_1(\tau_n^-), \ v_1'(\tau_n^+) = v_1'(\tau_n^-) + (-1)^{n+1}(\alpha_{n-1} + \alpha_n).$$
(18)

‡ We could also rescale the time variable $s = \alpha_0 \tau$ with $\tilde{v}_1(s) = v_1(\tau)$ to express (17) with only one parameter $\kappa = \mu/\alpha_0$, such that in the first time interval $\tilde{v}_1^2(s)\tilde{v}_1''(s) = -\kappa(\tilde{v}_1'(s) + \tilde{v}_1(s))$.

3.1. Solution in the time interval $0 \le \tau \le T$

In the first interval $0 \le \tau \le T$, we assume that the parameter $\kappa = \mu/\alpha_0$ is small. The initial conditions are given by $v_1(0) = 1$ and $v_1'(0) = \mu - \alpha_0 < 0$, as seen previously. We expect in this case that $v_1(\tau)$ decreases initially with $Q_1(\tau)$. The overall exact solution is given by an implicit relation between τ and v_1 . We first assume that v_1 can be expressed implicitly through the functional integral

$$\tau = \int_{v_1(\tau)}^1 \frac{v}{Z(v)} dv,\tag{19}$$

where Z(v) is an unknown function to be determined. By differentiating the previous expression with respect to τ , we have the identity $Z(v_1(\tau)) = -v_1(\tau)v_1'(\tau) > 0$. Expressing $v_1'(\tau)$ and $v_1''(\tau)$ in the equation (17) as functions of $Z(v_1(\tau))$ and $v_1(\tau)$, we obtain a new ODE of reduced order

$$Z'(v) - \frac{Z(v)}{v} = \mu \left(\frac{1}{v} - \frac{\alpha_0 v}{Z(v)}\right),\tag{20}$$

where Z'(v) is the derivative of Z with respect to v and not τ . The initial condition at $\tau = 0$ is replaced by $Z(1) = \alpha_0 - \mu > 0$, according to the previous values of $v_1(0)$ and $v_1'(0)$. We then perform the substitution $W(v) = (Z(v) + \mu)/v$ in order to simplify even further the previous ODE

$$W'(v) = \frac{\mu \alpha_0}{\mu - vW(v)}. (21)$$

This equation belongs to the class of Abel ODEs of the second kind [20, 21]. A way to resolve it is to use an inverse transformation, solving v as a function of W instead [22]. This yields the linear equation

$$\mu \alpha_0 v'(W) = \mu - W v(W), \tag{22}$$

whose solutions v(W) are simply given by the integrating factor

$$v \exp\left(\frac{W^2}{2\mu\alpha_0}\right) - \alpha_0^{-1} \int_0^W \exp\left(\frac{u^2}{2\mu\alpha_0}\right) du = C, \tag{23}$$

which determines W and Z as a function of v implicitly. C is a constant which can be evaluated by considering the identity $W(v_1(0)) = \alpha_0$ at $\tau = 0$. This gives the following value of C and its expansion for $\kappa = \mu/\alpha_0$ small

$$C = e^{1/2\kappa} - \int_0^1 \exp\left(\frac{u^2}{2\kappa}\right) du \simeq (1 - \kappa)e^{1/2\kappa}.$$
 (24)

A similar expansion for (23) can be performed when κ is small, assuming that μ is small compared to W

$$\left(v - \frac{\mu\alpha_0}{W(v)}\right) \exp\left(\frac{W^2(v)}{2\mu\alpha_0}\right) \simeq C.$$
 (25)

We can then obtain a further approximate value for Z(v), assuming that Z(v) always stays positive

$$Z(v) \simeq v\sqrt{2\mu\alpha_0\log\left(C/v\right)} - \mu. \tag{26}$$

This implies finally that the density $v_1(\tau)$ is given implicitly by

$$\tau \simeq \int_{v_1(\tau)}^1 \frac{v dv}{v \sqrt{-2\mu\alpha_0 \log\left(e^{-1/2\kappa}v\right)} - \mu}.$$
 (27)

In equation (19), we have assumed that $v_1(\tau) \leq 1$ is decreasing and that the function Z(v) has to be strictly positive, otherwise τ may not be defined if Z vanishes and becomes negative. For small values of τ , we can approximate (27) with $Z(v) \simeq \alpha_0 v$ and therefore obtain that $v_1(\tau) \simeq 1 - \alpha_0 \tau$, which follows the initial behavior of the function Q_1 and decreases as τ increases. However, for τ large, the inspection of equation (27) implies that $v_1(\tau)$ should reach asymptotically from above a limiting value $v^* > 0$ determined by the singularity of Z, and solution of $Z(v^*) = 0$:

$$v^* \sqrt{-2\mu\alpha_0 \log(e^{-1/2\kappa}v^*)} = \mu. \tag{28}$$

In the limit of κ small, we obtain that $v^* \simeq \kappa + \kappa^2 \log \kappa$. As τ increases, the integrand becomes indeed larger when v approaches v^* from above. This is equivalent to implying that $v_1(\tau) \to v^* > 0$ when $\tau \gg 1$, with T or α_0 large.

3.2. Numerical solution in the other intervals $\tau_n \leq \tau \leq \tau_{n+1}$

We study here the global solutions (16) for the successive time intervals, and for different sawtooth geometries such as the slope of each linear parts, T=1 being fixed. In regard to the time behavior of (9), we consider a diverging amplitude for the oscillating function $Q_1(\tau)$, by taking $\alpha_n = e^{\beta n}\alpha_0$, with $\beta > 0$. Figure 4 displays the solutions for different values of α_0 with $\beta = 0.1$. We have checked that using numerical algorithms for solving directly the non-local integral equation (10), or solving the ODE (16) with the Runge-Kutta-Fehlberg RKF45 method via Maple, leads to the same results. For small α_0 , v_1 follows the sawtooth behavior of Q_1 but stays positive. When α_0 increases, v_1 develops into a more complex periodic signal. It also takes very small values at times $\tau_k = (2k+1)T$, and its maximum amplitude is also smaller as additional frequency modes appear in the signal. It is more difficult to obtain solutions when $\alpha_0 > 2$, as v_1 takes very small values in time intervals that cannot be resolved precisely with the algorithms. In summary, v_1 is periodic and stays positive when we apply the approximate sawtooth kernel function Q_1 which is a diverging oscillating function with large negative and positive values. However additional modes appear in the time Fourier spectrum as α_0 increases, and which seem to be harmonics of the fundamental frequency π/T .

4. Model with a driven harmonic oscillator and negative damping parameter

In this section, we replace the kernels $Q_1(\tau)$ and $\tilde{Q}_1(\tau)$ by a function representing an oscillator characterized by a negative damping parameter. The idea is that to approximate $Q_1(\tau)$ and $\tilde{Q}_1(\tau)$ defined in equation (9) by a simple cosine function of frequency ω multiplied by a diverging exponential factor at rate γ . The quantities ω and γ can be viewed as the imaginary and real

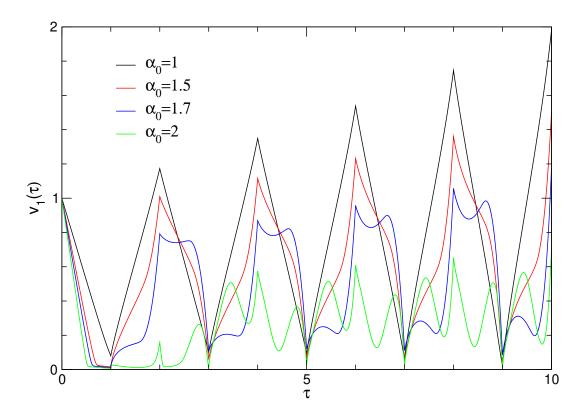


Figure 4: Solutions of the ODEs (16) for a regular sawtooth function Q_1 with $\alpha_n = \alpha_0 e^{\beta n}$, $\beta = 0.1$, $\mu = 0.04$, T = 1, and for different slope parameters α_0 .

parts of the largest eigenvalue z_{max}^+ . Although this approximation is correct asymptotically, it is not the case at small times since all the other eigenstates will give a finite contribution. However we expect to extract enough information about the physics of the system. We also introduce a phase shift $\varphi \in]-\pi/2,\pi/2[$ such that

$$Q_1(\tau) = e^{\gamma \tau} \frac{\cos(\omega \tau - \varphi)}{\cos \varphi}, \ \cos \varphi \neq 0, \tag{29}$$

with initial condition $Q_1(0) = 1$, independent of the phase. The phase diagram is actually symmetrical by the reflection operation $\varphi \to -\varphi$, so that we will restrict afterwards the analysis to the case when $\varphi \geq 0$. The derivative of this function at the origin is equal to $Q'_1(0) = \gamma + \omega \tan \varphi > 0$. Substituting (29) into equation (10), we obtain

$$v_1(\tau) = \frac{\mu}{\cos\varphi} \int_0^{\tau} \frac{e^{\gamma(\tau-\tau')}\cos(\omega(\tau-\tau')-\varphi)}{v_1(\tau')} d\tau' + e^{\gamma\tau} \frac{\cos(\omega\tau-\varphi)}{\cos\varphi}.$$
 (30)

We can use the fact that $Q_1'' - 2\gamma Q_1' + (\gamma^2 + \omega^2)Q_1 = 0$ to eliminate the dependence in Q_1 by considering the first and second derivative of the previous expression. After combining the different terms we finally obtain a second order ODE independent of the non-local integral

$$v_1''(\tau) - 2\gamma v_1'(\tau) + (\gamma^2 + \omega^2)v_1(\tau) = -\mu \frac{v_1'(\tau)}{v_1^2(\tau)} - (\gamma - \omega \tan \varphi) \frac{\mu}{v_1(\tau)}.$$
 (31)

The initial conditions are $v_1(0) = 1$ and $v'_1(0) = \gamma + \omega \tan \varphi + \mu$. Notice that we can remove the dependence in μ by substituting $v_1 \to \sqrt{\mu}v_1$. The left hand side is the linear ODE of an harmonic oscillator that $Q_1(\tau)$ satisfies, with frequency ω and negative damping parameter $-\gamma$, and the right hand side are non-linear terms. This equation has a time-independent solution in the long time limit provided that $\tan \varphi > \gamma/\omega$ and which is equal to

$$v_1^* = \left(\mu \frac{\omega \tan \varphi - \gamma}{\gamma^2 + \omega^2}\right)^{1/2}.$$
 (32)

In the opposite case, $\tan \varphi \leq \gamma/\omega$, the solution $v_1(\tau)$ decays to zero exponentially with τ , since for small v_1 the two terms on the right hand side of equation (31) are dominant, and we have instead

$$v_1'(\tau) \simeq -(\gamma - \omega \tan \varphi)v_1(\tau).$$
 (33)

This yields $v_1(\tau) \sim e^{-(\gamma - \omega \tan \varphi)\tau} \to 0$. To study the stability of the solution (32), we introduce a correction $w(\tau)$ to the steady state solution $v_1(\tau) = v_1^* + w(\tau)$, and linearize the equation (31)

$$w''(\tau) - \left(2\gamma - \frac{\mu}{v_1^{*2}}\right)w'(\tau) + 2(\gamma^2 + \omega^2)w(\tau) = 0.$$
 (34)

The solution w will grow if the damping factor is positive, or $2\gamma \ge \mu/v_1^{*2}$. This will occur if the phase φ is larger than a critical value φ_c which satisfies

$$\tan \varphi \ge \frac{3\gamma^2 + \omega^2}{2\gamma\omega} = \tan \varphi_c > \frac{\gamma}{\omega}.$$
 (35)

In this case, we observe steady oscillations with a frequency Ω determined approximately at this order by

$$\Omega^2 \simeq 2\gamma^2 + 2\omega^2,\tag{36}$$

and which is independent of the current μ as expected. In table 3 we have summarized the phase diagram as function of φ , according to the properties of (34).

Table 3

$0 \le \tan \varphi \le \gamma/\omega$	$v_1(\tau)$ goes to zero exponentially (33)	
$\gamma/\omega < \tan \varphi \le (3\gamma^2 + \omega^2)/2\gamma\omega$	$v_1(\tau)$ goes to the equilibrium value v_1^* (32)	
$(3\gamma^2 + \omega^2)/2\gamma\omega < \tan\varphi$	$v_1(\tau)$ oscillates with a finite amplitude and frequency Ω (36)	

In figure 5(a) is represented numerically the time evolution of v_1 as function of different values of the phase φ . There are two domains according to whether φ is lower or greater than φ_c ,

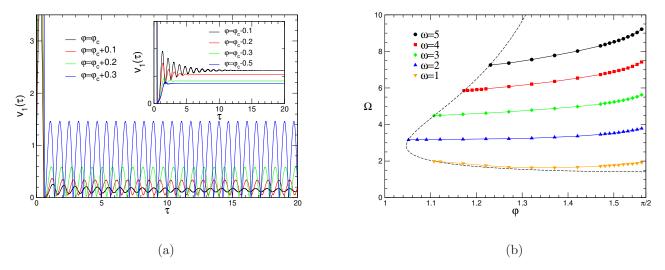


Figure 5: (a) Solution of the ODE (31) for different angles φ . The parameters are $\omega = 5$, $\gamma = 1$, and $\mu = 0.04$. The critical angle is $\varphi_c = \arctan(14/5) \simeq 1.228 > \arctan(\gamma/\omega) = \arctan(1/5)$. (b) Oscillation frequency Ω as function of φ for several values of ω , with parameters $\gamma = 1$, and $\mu = 0.04$. The dashed line is the threshold $\varphi = \varphi_c$ for which oscillations appear and which delimitates the domain of existence of Ω . This corresponds to $\Omega^2 = 2\gamma^2 + 2\omega^2$. In the limit $\varphi = \varphi_c = \pi/2$, $\omega \to 0$ and $\Omega = \sqrt{2}\gamma$.

for which v_1 reaches asymptotically the equilibrium value v_1^* or self-oscillates respectively. The oscillation frequency Ω can be extracted by Fourier transform of the time signals and represented as a curve function of φ for different values of ω , see figure 5(b). There is no simple relation between Ω and the other parameters, except at the critical angle φ_c for which $\Omega = \sqrt{2\gamma^2 + 2\omega^2}$, see equation (36).

Near φ_c , we can assume that the oscillation amplitude is small compare to the constant solution v_1^* . The previous linear expansion (34) shows that the damping factor $2\gamma - \mu/v_1^{*2}$ must be positive in order to generate oscillations. However, this expansion leads to an unstable solution as w diverges with time. We can study the next terms of the expansion and verify whether these terms stabilize the oscillation amplitude. If we expand (31) up to the third order in w, we obtain a non-linear equation with a time dependent damping factor

$$w''(\tau) - \left(2\gamma - \frac{\mu}{v_1^{*2}} + 2\frac{\mu}{v_1^{*3}}w(\tau) - 3\frac{\mu}{v_1^{*4}}w^2(\tau)\right)w'(\tau) + (2\gamma^2 + 2\omega^2)w(\tau) + (\gamma - \omega\tan\varphi)\frac{\mu}{v_1^{*3}}w^2(\tau) - (\gamma - \omega\tan\varphi)\frac{\mu}{v_1^{*4}}w^3(\tau) = 0.$$
 (37)

We have checked that this ODE displays stable oscillations close to the exact solution for $\varphi \gtrsim \varphi_c$, see figure 6(a). These oscillations are not symmetric, and their amplitude depends on ω , see figure

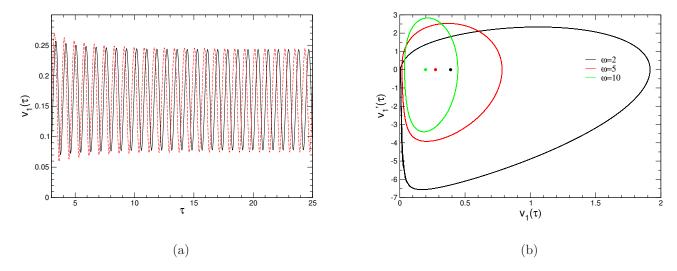


Figure 6: (a) Approximation corresponding to (37) (red dashed line) and comparison with the exact solution (31) (black line) for $\omega = 5$, $\gamma = 1$, and $\mu = 0.04$. $\varphi = \pi/2 - 0.3$ is chosen to be close to φ_c : $\delta \varphi = \varphi - \varphi_c \simeq 0.043$. (b) Oscillation cycles for three different parameters ω , with $\gamma = 1$, $\varphi = \pi/2 - 0.1$, and $\mu = 0.04$. The dot points correspond to the equilibrium solution v_1^* .

6(b). For example, as ω decreases, the amplitude increases. Restricting the expansion to the second order would still lead to non stable oscillations, due in particular to the presence of the linear term w(t) in the damping coefficient, in factor of w'(t). The stability comes instead from the quadratic term $w^2(t)$ in the same damping coefficient and from the last term in $w^3(t)$. To analyze the stability of (37), we consider an expansion in the small parameter $0 < \tan \varphi - \tan \varphi_c \simeq \delta \varphi \ll 1$. The technique we use in the following is described in the reference book [23] and can be applied to different ODEs with negative damping containing a linear or quadratic amplitude in addition to the constant term. This is usually useful to study self-excited oscillators. We first simplify the previous ODE in order to reduce the number of parameters. We chose to rescale the amplitude $w(\tau) = c\tilde{w}(\tau)$, with $c = v_1^* \omega \gamma \delta \varphi/(\gamma^2 + \omega^2)$, so that

$$\tilde{w}''(\tau) + \Omega_0^2 \tilde{w}(\tau) = 4\gamma \epsilon \left(1 + \tilde{w}(\tau) - \frac{3\epsilon}{2} \tilde{w}^2(\tau) \right) \tilde{w}'(\tau) + \frac{\epsilon}{2} \Omega_0^2 \left(\tilde{w}^2(\tau) - \epsilon \tilde{w}^3(\tau) \right), \quad (38)$$

with $\Omega_0^2 = 2(\gamma^2 + \omega^2)$, and $\epsilon = \omega \gamma \delta \varphi / (\gamma^2 + \omega^2) \ll 1$ the rescaled small parameter. This yields in particular $c = \epsilon v_1^*$. When $\epsilon = 0$, we can write the general solution (38) as $\tilde{w}(\tau) = u \cos(\Omega_0 \tau + \beta)$. When $\epsilon \neq 0$, the general solution can still be expressed in the same form, provided the coefficients u and β are time dependent. Since we have two functions satisfying one ODE, we can chose a constraint between u and β . It is usually convenient to impose the identity $\tilde{w}'(\tau) = -u\Omega_0 \sin(\Omega_0 \tau + \beta)$. This implies in particular that

$$u'(\tau)\cos\phi - u\beta'(\tau)\sin\phi = 0, (39)$$

where we used the notation $\phi = \Omega_0 \tau + \beta(\tau)$. Equation (38) can then be formally rewritten as

$$\tilde{w}''(\tau) + \Omega_0^2 \tilde{w}(\tau) = \epsilon f(\tilde{w}, \tilde{w}') = \epsilon f(u \cos \phi, -u\Omega_0 \sin \phi), \tag{40}$$

with function f defined by

$$f(\tilde{w}, \tilde{w}') = 4\gamma \left(1 + \tilde{w} - \frac{3\epsilon}{2}\tilde{w}^2\right)\tilde{w}' + \frac{1}{2}\Omega_0^2 \left(\tilde{w}^2 - \epsilon\tilde{w}^3\right). \tag{41}$$

After computing the second derivative $w''(\tau)$ from the expression $w'(\tau) = -u\Omega_0 \sin \phi$, we obtain

$$\tilde{w}''(\tau) + \Omega_0 \tilde{w}(\tau) = -u'(\tau)\Omega_0 \sin \phi - u\Omega_0 \beta'(\tau) \cos \phi. \tag{42}$$

The right hand side can be identified with (40), and using the relation between u and β , see (39), we obtain the following set of first order ODEs

$$u' = -\frac{\epsilon}{\Omega_0} \sin \phi f(u \cos \phi, -u\Omega_0 \sin \phi),$$

$$\beta' = -\frac{\epsilon}{\Omega_0 u} \cos \phi f(u \cos \phi, -u\Omega_0 \sin \phi).$$
(43)

These equations are identical to (38), since no approximation have been made yet. We then assume that ϕ varies rapidly compared to the variations of u and β . We can replace in this approximation the previous equations by their time average over a period $2\pi/\Omega_0$, after considering the quantities u' and β' quasi constant. This yields

$$u' \simeq -\frac{\epsilon}{2\pi\Omega_0} \int_0^{2\pi} d\phi \sin\phi f(u\cos\phi, -u\Omega_0\sin\phi),$$

$$\beta' \simeq -\frac{\epsilon}{2\pi\Omega_0 u} \int_0^{2\pi} d\phi \cos\phi f(u\cos\phi, -u\Omega_0\sin\phi).$$
 (44)

After performing the integrations over ϕ , we obtain a set of coupled ODEs for the amplitude and phase shift

$$u' = \gamma \epsilon (2u - \frac{3\epsilon}{4}u^3), \ \beta' = \frac{3}{16}\Omega_0 \epsilon^2 u^2. \tag{45}$$

We should notice that the evolution of the amplitude u does not depend on the terms in factor of Ω_0^2 in equation (41), after averaging over ϕ , but is determined by the damping terms in factor of \tilde{w}' , especially the constant and quadratic term \tilde{w}^2 . However, the change in frequency depends instead on the cubic term in factor of Ω_0^2 . All the terms which are a odd function of ϕ do not contribute to (45). The presence of the quadratic term in the damping factor leads to the existence of a finite equilibrium limit for u which converges exponentially to a positive constant u_0 determined by the non-zero solution of u' = 0 in (45), or explicitly $u_0 = \sqrt{8\ell/3\epsilon}$. This leads, after rescaling to the original variable, to the equilibrium amplitude $w_0 = \sqrt{8\epsilon/3}v_1^*$, and therefore we can conclude that the effect of the additional quadratic term in the damping factor is to stabilize the oscillations that are generated self-consistently. The frequency Ω is modified accordingly after solving (45) for β' in the asymptotic limit, and is shifted by a small amount such that the corrected frequency is given by $\Omega \simeq \Omega_0 + \beta'$ or $\Omega \simeq \Omega_0(1 + \epsilon/2)$. It is, as expected, independent of μ , contrary to the amplitude which is proportional to $v_1^* \propto \sqrt{\mu}$. As an example, if we take the parameters from

figure 5(b), $\omega = 3$, $\gamma = 1$ and $\mu = 0.04$, we obtain for $\varphi = 1.17 \gtrsim \varphi_c \simeq 1.107$ the overall estimated amplitude of the signal $2w_0 \simeq 0.185$, which is close to the numerical value extracted from the analysis of the ODE (38) or 0.217, with $\epsilon = 0.128$. We also obtain $\Omega \simeq 4.759$, to be compare with $\Omega_0 \simeq 4.472$ and the numerical estimation from the Fourier transform of the signal $\Omega \simeq 4.634$.

The present model in this section does not consider all the eigenstates with eigenvalues that have negative real part, since we have assumed that these states have a negligible contribution as they correspond to the terms in (9) that vanish exponentially with time. However their contribution at small times may be relevant to the general solution as the convolution integral equations (8) and (10) are non local. In the following section, we will generalize the previous method to the main model which incorporates N oscillators.

5. Generalization to a set of harmonic oscillators

The generalization of the previous results to the main model of section 2 is important in order to analyze the onset of the self-oscillations, and to extract the value of the frequency observed for example in figure 3 and which can not be approximated with a simple oscillator. We first write a general ODE for v_1 by removing all dependence in the Q_1 and \tilde{Q}_1 functions, as well as the convolution integrals. We expect that for N oscillators, v_1 satisfies a N^{th} -order ODE. As in the previous section, we can write an ODE satisfied by Q_1 and \tilde{Q}_1 , since they are a sum of N harmonic oscillators (or N-1 precisely because of the presence of the zero eigenvalue eigenstate). From equation (9), it is easy to establish a linear ODE satisfied by Q_1 and \tilde{Q}_1 using the characteristic polynomial of the coefficient matrix A whose coefficients are given by

$$\det\left(zI_N - A\right) = \sum_{k=0}^{N} a_k z^k,\tag{46}$$

with $a_N = 1$ and $a_0 = 0$. Since the eigenvalues are roots of this polynomial, we have directly an ODE for Q_1 and \tilde{Q}_1 after differentiating (9)

$$\sum_{k\geq 1}^{N} a_k Q_1^{(k)}(\tau) = 0. (47)$$

By differentiating N times the identity (10), and multiplying each successive equations by a_k , we can eliminate all the time dependence involving Q_1 and \tilde{Q}_1 , and therefore all the convolution integrals disappear. This yields

$$\sum_{k=1}^{N} a_k \left(v_1^{(k)} - \mu \sum_{l=0}^{k-1} Q_1^{(k-1-l)}(0) \frac{\partial^l}{\partial \tau^l} \left(\frac{1}{v_1} \right) \right) = 0.$$
 (48)

This equation contains only derivatives in v_1 since $a_0 = 0$ and also because there is no contribution in v_1 coming from the l = 0 terms in the previous sum as

$$\sum_{k=1}^{N} a_k Q_1^{(k-1)}(0) = 0. (49)$$

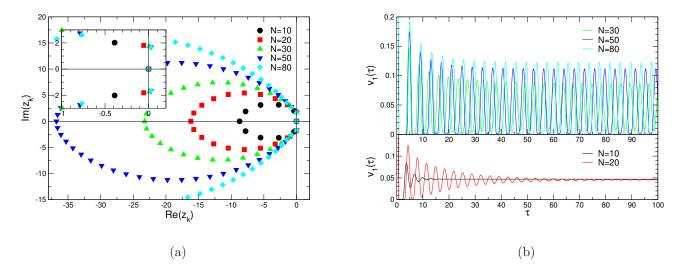


Figure 7: (a) Plot of the complex roots of R(z), equation (54), for different values of N. The parameters are $\mu = 2.10^{-4}$, $\lambda = 0.06$ and a = 0.9. In inset is represented the roots near the origin. Only the cases N = 30, N = 50, and N = 80 have a pair of roots with strictly positive real parts. (b) Time dependence of $v_1(\tau)$ for the previous cases. Oscillations are observed for N = 30, 50, 80. For the other cases $N = 10, 20, v_1$ reaches the asymptotic limits $v_1^* \simeq 0.047$ and $v_1^* \simeq 0.045$ respectively, see equation (12).

This comes from the fact that one of the factors α_l in the sum (9), $Q_1(\tau) = \sum_{l=1}^N \alpha_l e^{z_l \tau}$, vanishes when $l = l_0$, with l_0 associated to the zero eigenvalue $z_{l_0} = 0$. This has been checked numerically. Otherwise, it is easy to see that the previous equation (48) would not be consistent with the existence of the time-independent solution that we found in section 2.2. There is therefore no v_1 term in the ODE, contrary to the case of the single harmonic oscillator that we studied in the previous section. Each derivative term of the ODE (48) can be rearranged by increasing differential order of v_1^{-1}

$$\sum_{k=1}^{N} a_k v_1^{(k)} - \mu \left(\frac{1}{v_1}\right)' \left(a_2 Q_1(0) + a_3 Q_1'(0) + \dots + a_N Q_1^{(N-2)}(0)\right)$$

$$-\mu \left(\frac{1}{v_1}\right)'' \left(a_3 Q_1(0) + a_3 Q_1'(0) + \dots + a_N Q_1^{(N-3)}(0)\right) - \dots$$

$$-\mu \left(\frac{1}{v_1}\right)^{(N-1)} a_N Q_1(0) = 0. \tag{50}$$

We can define (N-1) additional coefficients $b_k = \sum_{l=0}^{N-k-1} a_{k+1+l} Q_1^{(l)}(0)$ such that the previous ODE becomes

$$\sum_{k=1}^{N} a_k v_1^{(k)} - \mu \sum_{k=1}^{N-1} b_k \left(\frac{1}{v_1}\right)^{(k)} = 0.$$
 (51)

We can, as before, eventually remove the dependence on μ by rescaling $v_1 \to \sqrt{\mu}v_1$. The complex coefficients b_k can be computed recursively as the inverse powers of the eigenvalues, using the factors α_l . For the first two terms, we have indeed

$$b_1 = -a_1 \sum_{l \neq l_0} \alpha_l z_l^{-1}, \ b_2 = -a_1 \sum_{l \neq l_0} \alpha_l z_l^{-2} - a_2 \sum_{l \neq l_0} \alpha_l z_l^{-1}, \tag{52}$$

and for the general case $b_k = -\sum_{i=1}^k a_i \sum_{l \neq l_0} \alpha_l z_l^{-k-1+i}$. From these results, we can analyze the stability of the steady state solution v_1^* of the equation (12) if we linearize (51) by setting $v_1 = v_1^* + w$, $|w| \ll v_1^*$, which yields

$$\sum_{k=1}^{N} a_k w^{(k)} + \frac{\mu}{v_1^{*2}} \sum_{k=1}^{N-1} b_k w^{(k)} = 0.$$
 (53)

The stability of the solution can be studied by finding the roots of the following polynomial

$$R(z) = \sum_{k=1}^{N-1} \left(a_k + \frac{\mu}{v_1^{*2}} b_k \right) z^k + a_N z^N.$$
 (54)

This polynomial is different from the characteristic polynomial of A due to the additional terms b_k . The perturbation $w(\tau)$ will grow if at least one complex root of R(z) has a strictly positive real part. In figure 7(a) we have represented the roots of R(z) for different N, with a fixed set of parameters. As N increases, one pair of conjugate roots moves into the positive part of the real axis, and self-oscillations occur, see figure 7(b). The imaginary part of the root z_{max}^+ of R(z) with the largest real part gives an estimate of the oscillation frequency. We find for N=50 $z_{max}^+ \simeq 0.092 + 1.454i$, and the Fourier transform of the signal of figure 7(b) gives approximately $\Omega \simeq 1.332$ for this value of N. Im (z_{max}^+) is closer to Ω than the value given by the characteristic polynomial of A which is equal to Im $(z_{max}^+) \simeq 1.561$.

A more precise analysis of (51) requires to estimate the contributions from all the non-linear terms. The successive derivatives of v_1^{-1} can be ordered as function of the powers of v_1^{*-1} . This can be done recursively by using a combinatorial method

$$\sum_{k=1}^{N} a_k w^{(k)} + \mu \sum_{k=1}^{N-1} b_k \left(\frac{1}{v_1^{*2}} \frac{w^{(k)}}{(1+w/v_1^*)^2} - \frac{1}{v_1^{*3}} \sum_{l_1=1}^{k-1} \binom{k}{l_1} \frac{w^{(l_1)} w^{(k-l_1)}}{(1+w/v_1^*)^3} \right)$$
(55)

$$+\frac{1}{v_1^{*4}} \sum_{l_1=1}^{k-1} \sum_{l_2=1}^{k-l_1-1} \binom{k}{l_1} \binom{k-l_1}{l_2} \frac{w^{(l_1)}w^{(l_2)}w^{(k-l_1-l_2)}}{(1+w/v_1^*)^4} - \dots + \frac{(-1)^{k+1}}{v_1^{*k+1}} \frac{k!w^{(1)^k}}{(1+w/v_1^*)^{k+1}} = 0.$$

This expression is exact. This gives a formal expansion in v_1^{*-1} . When v_1^* is large enough, we can approximate this ODE by retaining only the first few terms. As in section 4, we can analyze the

stability of the solution associated to the pair of conjugate roots $z_{max}^{\pm} = \Lambda_0 \pm i\Omega_0$ of the polynomial (54) with Λ_0 assumed in the following to be strictly positive but smaller than Ω_0 . We consider the perturbation amplitude $\tilde{w}(\tau) = e^{\Lambda} \cos(\Omega_0 \tau + \beta)$, after rescaling $w = \sqrt{\mu}\tilde{w}$ and set $\epsilon = \sqrt{\mu}/v_1^*$ which is the small parameter of the expansion. The quantities $\Lambda = \Lambda_0 \tau + a(\tau)$ and β depend on time and a perturbative solution can be given as an expansion at the lowest order in ϵ . We also choose a constraint between the two functions Λ and β and impose the following identity for the derivative: $\tilde{w}'(\tau) = e^{\Lambda}(\Lambda_0 \cos \phi - \Omega_0 \sin \phi)$, where we used the notation $\phi = \Omega_0 \tau + \beta(\tau)$. This implies the following condition

$$a'(\tau)\cos\phi - \beta'(\tau)\sin\phi = 0. \tag{56}$$

After some calculations and keeping only the linear terms in a' and β' , we establish the following successive derivatives of \tilde{w}

$$\tilde{w}^{(k)}(\tau) \simeq e^{\Lambda} \left[Z_k(\phi) + c_k \left(a'(\tau) \sin \phi + \beta'(\tau) \cos \phi \right) \right],$$

$$Z_k(\phi) = \text{Re} \left[(z_{max}^+)^k e^{i\phi} \right], \ c_k = -\text{Im} \left(z_{max}^+ \frac{(z_{max}^+)^{k-1} - \Lambda_0^{k-1}}{z_{max}^+ - \Lambda_0} \right).$$
(57)

We then replace the derivatives of \tilde{w} in the ODE (55) by their expression from (57), and organize the different terms as a series expansion in ϵ . We obtain a single equation for the combination $a'(\tau) \sin \phi + \beta'(\tau) \cos \phi$:

$$(a'(\tau)\sin\phi + \beta'(\tau)\cos\phi) \sum_{k=1}^{N-1} c_k (a_k + \epsilon^2 b_k + \mathcal{O}(\epsilon^3)) \simeq$$

$$\epsilon^3 e^{\Lambda} \sum_{k=1}^{N-1} b_k \sum_{l=1}^{k-1} \left(\binom{k}{l} Z_l Z_{k-l} + 2Z_k \cos\phi \right)$$

$$-\epsilon^4 e^{2\Lambda} \sum_{k=1}^{N-1} b_k \left(\sum_{l_1=1}^{k-1} \sum_{l_2=1}^{k-l_1-1} \binom{k}{l_1} \binom{k-l_1}{l_2} Z_{l_1} Z_{l_2} Z_{k-l_1-l_2} \right)$$

$$+3 \sum_{l=1}^{k-1} \binom{k}{l} Z_l Z_{k-l} \cos\phi + 3Z_k \cos^2\phi + \mathcal{O}(\epsilon^5). \tag{58}$$

As before, we can multiply this expression by $\sin \phi$ and $\cos \phi$, and use the relation (56) to obtain independently an equation for a' and β' respectively. A time average over the period $2\pi/\Omega_0$ is then performed, after we assumed that a' and β' are quasi constant quantities. After some algebra, we obtain the set of ODEs

$$a'(\tau) \sum_{k=1}^{N-1} c_k (a_k + \epsilon^2 b_k) \simeq -\frac{3\epsilon^4}{8} e^{2\Lambda} \sum_{k=1}^{N/2} (-1)^k b_{2k-1} \Omega_0^{2k-1}$$

$$\beta'(\tau) \sum_{k=1}^{N-1} c_k (a_k + \epsilon^2 b_k) \simeq -\frac{3\epsilon^4}{8} e^{2\Lambda} \sum_{k=1}^{(N-1)/2} (-1)^k b_{2k} \Omega_0^{2k}, \tag{59}$$

where we have dropped the dependence in $\Lambda_0 \ll \Omega_0$ in the coefficients in order to simplify the expressions. We can rewrite the first equation using the substitution $u(\tau) = e^{\Lambda}$, and obtain the

relation $u' \simeq \Lambda_0 u - 3\epsilon^3 a_0 u^3/8$, with the constant a_0 given by

$$a_0 = \frac{\sum_{k=1}^{N/2} (-1)^k b_{2k-1} \Omega_0^{2k-1}}{\sum_{k=1}^{N-1} c_k (a_k + \epsilon^2 b_k)}.$$
(60)

u reaches an asymptotic non-zero limit if $a_0 > 0$, given by $u_0^2 = 8\Lambda_0/(3a_0\epsilon^4)$. If $\Lambda_0 < 0$, $u_0 = 0$ and no oscillation is observed. If $\Lambda_0 = 0$ and $a_0 < 0$, additional corrections are necessary in order to verify the existence of oscillations. The equation for the frequency shift can be written as $\beta' \simeq -3\epsilon^4 u^2 \beta_0/8$, with constant β_0 given by

$$\beta_0 = \frac{\sum_{k=1}^{(N-1)/2} (-1)^k b_{2k} \Omega_0^{2k}}{\sum_{k=1}^{N-1} c_k (a_k + \epsilon^2 b_k)}.$$
(61)

The shift in frequency will therefore be given asymptotically by $\beta' = -\Lambda_0 \beta_0/a_0$, independent of ϵ . The sign of the shift depends on the sign of β_0 , in the case where $a_0 > 0$. We can extend the expansion (58) up to the ϵ^6 order, and find that $\beta' \simeq -\epsilon^4 u^2 \beta_0 (3/8 + 5\epsilon^2 u^2/16)$, which should give a negligible correction. For example, for the parameters given in figure 7 and N = 80, we obtain numerically $\epsilon \simeq 0.319$, which is small enough to apply the previous expansion. In particular, after solving for the roots of the polynom R, we find $z_{max}^+ = \Lambda_0 + i\Omega_0 \simeq 0.026 + 1.669i$, and therefore $\Lambda_0/\Omega_0 \ll 1$. The frequency Ω from the numerical resolution of the ODEs (2) was evaluated by Fourier analysis and found to be equal to $\Omega \simeq 1.562$, while the correction $\Omega_0 + \beta' \simeq 1.517$ gives a closer value to the numerical estimate than Ω_0 . In this case we find that $a_0 > 0$ and $\beta_0 > 0$. The amplitude of the oscillations, see figure 7(b), corresponding to $\sqrt{\mu}u_0$, was found after a non-linear curve fitting with a simple cosine function to be equal approximately to 0.071 while the amplitude given by the expansion (59) gives 0.047, for $v_1^* \simeq 0.044$, which is of the same order of magnitude. The expansion (59) gives a first estimate of the oscillation characteristics, and shows that these oscillations can be stabilized by considering the first non-linear terms of the ODE (51).

6. Conclusion

In this paper, we have considered a simple model of coagulation-fragmentation of particle clusters which displays collective oscillations for a distribution of masses broad enough. The mass-dependent kernel leads to a coefficient matrix that has at least two conjugate complex eigenvalues with strictly positive real part, provided the size of the matrix is sufficiently large. This is a necessary condition for the emergence of a self-oscillating mechanism. This model allows us to evaluate all the densities from the monomer density solution v_1 that satisfies an implicit non-local Volterra integral equation. This equation depends on two kernels or excitation functions (which are identical in the large N limit) that have a strongly exponential divergence with positive and negative values. However v_1 appears to oscillate with a finite amplitude and with a minimum positive value proportional to the current. The oscillation amplitude depends on the parameter values such as current, the exponent of the coagulation kernel and the fragmentation rate. When

§ The numerical coefficients 3/8 and 5/16 are the angle average of $\cos^4 \phi$ and $\cos^6 \phi$ respectively.

we approximate the kernels with a linear or sawtooth function, we have shown that v_1 can never be negative, and is larger than μ when μ is small in the first approximation.

When we consider a harmonic function as an approximation of the kernels in the long time limit, i.e. with a diverging oscillatory behavior, the model shows self-excited oscillations that can be explained by the presence of a fluctuating negative damping in the expansion of the corresponding ODE around the threshold beyond which oscillations appear. This phenomenon is very similar to the behavior of a van der Pol oscillator in non-conservative systems which has a quadratic contribution in the damping factor. When the amplitude increases, the damping factor tends to be positive because of the minus sign in front of the quadratic term, and therefore the amplitude is reduced. And inversely, when the amplitude is small, the damping factor is negative and tends to increase the overall amplitude. The non-conservative oscillator can sustain self-oscillations only in the presence of a current, which drives the amplitude. The period of these oscillations does not depend on the current, however no exact solution can be found for its dependence on the other parameters, except around the oscillation threshold using a perturbative calculus, where it is found to be equal in this approximate model to $\sqrt{2}$ times the modulus of the largest eigenvalue of the coefficient matrix. The generalization to N oscillators leads to a more complex ODE that can be solved using the same perturbative method for some range of small parameters, but it is also the presence of negative damping factors that gives rise to self-oscillations.

Several questions remain concerning the presence of self-oscillations observed in more general aggregation kernels and without source [12]. In the present manuscript, the oscillations are driven by a source of monomers, but a more comprehensive study is needed for closed and conservative systems yielding oscillations without any external source, including extended fragmentation processes between clusters of different masses.

Acknowledgments

I would like to thank the Department of Physics and Astronomy at the Seoul National University where part of this work was carried out for its hospitality.

References

- [1] Krapivsky P L, Redner S and Ben-Naim E 2010 A Kinetic View of Statistical Physics (New York: Cambridge University Press)
- [2] Cuzzi J N, Burns J A, Charnoz S, Clark R N, Colwell J E, Dones L, Esposito L W, Filacchione G, French R G, Hedman M M, Kempf S, Marouf E A, Murray C D, Nicholson P D, Porco C C, Schmidt J, Showalter M R, Spilker L J, Spitale J N, Srama R, Sremčević M, Tiscareno M S and Weiss J 2010 Science 327 1470–1475
- [3] Brilliantov N, Krapivsky P L, Bodrova A, Spahn F, Hayakawa H, Stadnichuk V and Schmidt J 2015 Proceedings of the National Academy of Sciences 112 9536–9541
- [4] Brilliantov N V, Formella A and Pöschel T 2018 Nat. Commun. 9 797
- [5] White H W 1981 J. Colloid Interface Sci. 87 204–208
- [6] Ziff R M, McGrady E D and Meakin P 1985 J. Chem. Phys. 82 5269-5274

- [7] da Costa F P 2015 Mathematical Aspects of Coagulation-Fragmentation Equations *Math. Energy Clim. Chang.* ed Bourguignon J P, Jeltsch R, Pinto A and Viana M (Springer) pp 83–162 cim series ed
- [8] Krapivsky P L, Otieno W and Brilliantov N V 2017 Phys. Rev. E 96 042138
- [9] Bodrova A S, Stadnichuk V, Krapivsky P L, Schmidt J and Brilliantov N V 2019 J. Phys. A Math. Theor. 52 205001
- [10] Ball R C, Connaughton C, Jones P P, Rajesh R and Zaboronski O 2012 Phys. Rev. Lett. 109 168304
- [11] Connaughton C, Dutta A, Rajesh R, Siddharth N and Zaboronski O 2018 Phys. Rev. E 97 022137 (Preprint 1710.01875)
- [12] Brilliantov N V, Otieno W, Matveev S A, Smirnov A P, Tyrtyshnikov E E and Krapivsky P L 2018 *Phys. Rev.* E 98 012109
- [13] Buchstaber V M and Bunkova E Y 2012 Algebraically integrable quadratic dynamical systems (*Preprint* 1212.6675) URL http://arxiv.org/abs/1212.6675
- [14] Calogero F, Conte R and Leyvraz F 2020 J. Math. Phys. 61 102704
- [15] Wattis J A 2006 Physica D $\mathbf{222}$ 1–20
- [16] Vigil R D, Ziff R M and Lu B 1988 Phys. Rev. B 38 942
- [17] Hayakawa H 1987 J. Phys. A 20 L801–L805
- [18] Hayakawa H and Hayakawa S 1988 Publ. Astron. Soc. Japan 40 341–345
- [19] Burton T A 2005 Volterra Integral and Differential Equations 2nd ed (Mathematics in Science and Engineering vol 202) (Amsterdam: Elsevier)
- [20] Polyanin A D and Zaitsev V F 2003 Handbook of exact solutions for ordinary differential equations (Chapman & Hall/CRC)
- [21] Murphy G M 1960 Ordinary Differential Equations and Their Solutions (New York: Van Nostrand Reinhold Company)
- [22] Cheb-Terrab E S and Roche A D 2003 Eur. J. Appl. Math. 14 217-229
- [23] Nayfeh A H and Mook D T 1995 Nonlinear oscillations Wiley Classics Library (New York: Wiley) pages 121 and 129